

CRITICAL MEASURES FOR VECTOR ENERGY: GLOBAL STRUCTURE OF TRAJECTORIES OF QUADRATIC DIFFERENTIALS

ANDREI MARTÍNEZ-FINKELSHTEIN AND GUILHERME L. F. SILVA

ABSTRACT. Saddle points of a vector logarithmic energy with a vector polynomial external field on the plane constitute the *vector-valued critical measures*, a notion that finds a natural motivation in several branches of analysis. We study in depth the case of measures $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$ when the mutual interaction comprises both attracting and repelling forces.

For arbitrary vector polynomial external fields we establish general structural results about critical measures, such as their characterization in terms of an algebraic equation solved by an appropriate combination of their Cauchy transforms, and the symmetry properties (or the S -properties) exhibited by such measures. In consequence, we conclude that vector-valued critical measures are supported on a finite number of analytic arcs, that are trajectories of a quadratic differential globally defined on a three-sheeted Riemann surface. The complete description of the so-called critical graph for such a differential is the key to the construction of the critical measures.

We illustrate these connections studying in depth a one-parameter family of critical measures under the action of a cubic external field. This choice is motivated by the asymptotic analysis of a family of (non-hermitian) multiple orthogonal polynomials, that is subject of a forthcoming paper. Here we compute explicitly the Riemann surface and the corresponding quadratic differential, and analyze the dynamics of its critical graph as a function of the parameter, giving a detailed description of the occurring phase transitions. When projected back to the complex plane, this construction gives us the complete family of vector-valued critical measures, that in this context turn out to be vector-valued equilibrium measures.

1. INTRODUCTION AND STATEMENT OF MAIN RESULTS

1.1. Historical background. Statistical systems of many particles have been object of intense analysis for a long time, both from the perspective of physics and mathematics. The study of a particular type of interacting particle systems, the so-called determinantal point processes (and their cousins, Pfaffian processes), has been especially fruitful in the past thirty years. This success can be explained both by the ubiquitous character and flexibility of these models (describing the eigenvalues of several random matrix ensembles, non-intersecting diffusion paths, random growth models, random tilings, to mention a few) and the introduction of new tools,

2010 *Mathematics Subject Classification.* Primary: 31A15; Secondary: 14H05, 30C70, 30E10, 30F10, 30F30.

Key words and phrases. Logarithmic potential theory, vector energy, extremal problems, equilibrium on the complex plane, critical measures, S -property, quadratic differentials trajectories, compact Riemann surface.

intrinsically related with the analytic theory of orthogonal polynomials and their generalizations.

A common feature of these models is that either the joint probability density, the correlation functions, the normalization constant or the generating function can be expressed as a determinant (resp., a Pfaffian), and the right selection of the functions appearing in these determinants unveils the integrable structure of the underlying processes. Another unified property of these models is the possibility to put them in the framework of the so-called log gases, where the particles behave as charges on one or two dimensional sets, subject to the logarithmic interaction.

The best known case is the spectrum of some unitary matrix ensembles [28, 29], described in terms of classical families of orthogonal polynomials on the real line. Their zeros, all real and simple, satisfy an electrostatic model that goes all the way back to Stieltjes [76], and solve a minimization problem for the associated (logarithmic) energy. In other words, the zeros provide an equilibrium configuration on the conducting real line. Similar situation occurs for some non-intersecting paths models or for the six-vertex model in statistical mechanics [24], to mention some more examples.

Another classical framework for orthogonal polynomials (orthogonality on the unit circle), developed in the seminal works of Szegő [78], is connected to the analysis of the Ising model [30, 31].

Further immediate generalizations of the problems above oblige us to leave the real line and extend the notion of orthogonality and the associated log gas models to the complex plane. The so-called *non-hermitian orthogonal polynomials* appear naturally in approximation theory, as denominators of rational (e.g. Padé) approximants to analytic functions [40, 63], or equivalently, in the study of continued fractions. Electrostatic models on a conducting plane are also associated with the polynomial solutions of certain linear differential equations [58, 60].

Recently, non-hermitian orthogonality found its way to areas with a much more “modern” flavor, playing the crucial role in the description of the rational solutions to Painlevé equations [12, 17], in theoretical physics [1, 2, 20, 25] and in numerical analysis [27].

Observe that due to analyticity, there is a certain freedom in the choice of the integration contour for the non-hermitian orthogonal polynomials, which means that the location of their zeros is a priori not clear. As it follows from the fundamental works of Nuttall [64] and Stahl [74, 75] (see also the survey [4]), the limiting zero counting measure of these polynomials can be characterized in terms of certain max-min problems of the logarithmic potential theory, where the energy is minimized on each compact set, but these minima are maximized over an appropriately chosen class of contours. These questions (whose simplest example is the Chebotarev’s problem [55, 65] about the selection of sets of minimal logarithmic capacity) are closely related to the free boundary problems in the theory of partial differential equations.

Motivated by his study of convergence of Padé approximants to functions with branch points, Stahl [74, 75] not only proved that the zeros of the associated non-hermitian orthogonal polynomials distribute asymptotically as the equilibrium measure on the set of minimal capacity within the appropriate class (proving in passing their existence), but also characterized these sets in terms of a symmetry property, now known as the *S-property*. These ideas were considerably extended by Gonchar

and Rakhmanov [40] to include the case when the orthogonality weight depends on the degree of the polynomial. The electrostatic model features now a non-trivial external field (or background potential) that persists in the asymptotic regime; in consequence, the notion of the S -property had to be extended to deal with this external field.

The resulting Gonchar-Rakhmanov-Stahl (or GRS) theory allows to formulate conditional theorems of the form: “if there exists a compact set with the S -property corresponding to the associated electrostatic model, then the weak limit of the zero counting measures of the related non-hermitian orthogonal polynomials is the equilibrium measure of this set”. However, the existence of such compact sets is a problem from the geometric function theory, and it usually has to be settled in each concrete situation, basically from scratch. There are very few genuinely general rigorous results in this sense, for instance the most recent one of Kuijlaars and Silva [53] settling down the existence for polynomial external fields on the plane, and also the contribution by Kamvissis and Rakhmanov [45, 46] which deals with the Green potential instead of the logarithmic potential.

The notion of the S -property found also a natural interpretation in the light of the Deift-Zhou’s non-linear steepest descent method for the Riemann-Hilbert problems [28]. One of the key steps in the asymptotic analysis is the deformation of the contours, which transforms oscillatory behavior into a non-oscillatory plus exponentially decaying one. This deformation consists in “opening of lenses” along the level sets of certain functions (known as g -functions). It is precisely the S -property of these sets which guarantees the exponential decay on all non-relevant contours.

The available log gas models on the conducting plane suggest to extend the analysis of the equilibrium distributions (local energy minima on the prescribed sets) to the study of saddle points of the logarithmic energy on the plane. This yields the notion of *critical measures*, studied systematically in [58], following the original ideas in the unpublished manuscript [66]. The advantage of dealing with saddle points of the energy is the possibility to avoid the complications of the free boundaries, as well as to use variational techniques that go back to the work of Löwner and Schiffer.

One of the main conclusions of [58] (see also [53]) is the fact that the Cauchy transform of any critical measure satisfies (a.e. on the plane) a quadratic equation, known as the *spectral curve* in the context of random matrix theory. An immediate consequence of this fact is that the support of such a measure is a finite union of analytic curves, all of them being trajectories of the same quadratic differential on the plane.

Quadratic differentials are known to play an essential role in the study of conformal and quasi-conformal mappings, univalent functions, etcetera. Teichmüller [79], based on the fundamental work of Grötzsch, formulated a principle according to which quadratic differentials and their singularities are associated with the solutions of extremal problems in geometric function theory. Although he did not prove any general result supporting his principle, we should give credit to Teichmüller for putting spotlight on quadratic differentials. Schiffer [71, 72] introduced several variational techniques in the geometric function theory, allowing to find differential equations for the boundaries of the extremal domains, which in turn give a more formal connection with quadratic differentials. See [43, 56, 67] for further details.

As it was mentioned above, the Riemann-Hilbert asymptotic analysis for 2×2 matrix-valued functions characterizing the non-hermitian orthogonality provides another natural connection with quadratic differentials: the “right” level curves for the g -functions, that usually have a quite complicated structure, are trajectories of a suitable quadratic differential. Thus, the problem of existence of the appropriate curves can be reduced to the question about the global structure of such trajectories. This is not a simplification (the description of the global structure can be a formidable task), but it opens the gate to other techniques from the geometric function theory. These ideas have been successfully applied by Bertola [15], being closely related (in fact, in some sense equivalent) to the max-min approach mentioned before, see also [18]. In the same spirit and in a much more explicit form, quadratic differentials also appeared in similar asymptotic problems in [9, 10, 42, 49, 50, 57]. We emphasize that in all these problems, the underlying quadratic differentials are defined on the complex plane.

Returning to our motivation rooted in the interacting particle systems, recent works have showed that more general models, such as the hermitian matrix model with an external source [21–23], the two-matrix model [11, 16, 32–34], some classes of non-intersecting paths [47, 48] or the normal matrix model [26, 52] require considering log gases with more than one group of charges. In this case, standard orthogonal polynomials are not enough, and we must turn to the so-called *multiple orthogonal polynomials* and to the associated vector equilibrium problems [39]. The connection of these two notions was put forward first by Gonchar and Rakhmanov [38, 39], motivated by the analytic theory of Hermite–Padé approximants. The reader interested in further details is advised to look into the monograph by Nikishin and Sorokin [63], which features a nice introduction to vector equilibrium problems from the approximation theory’s perspective.

Although the analytic theory of multiple orthogonal polynomial is still in its infancy, two somewhat extreme cases are fairly well understood (see e.g. [5, 80]). The first one, known as the Angelesco system, corresponds to orthogonality with respect to measures living on disjoint sets, whose asymptotics is described by a vector equilibrium for repelling systems of charges. The other one, called the Nikishin system, where on the contrary the orthogonality is defined on the same set for all measures, is governed by a system of mutually attracting charges. In both situations the size of the charged ensembles is fixed. From the recent works of Aptekarev [3] and Rakhmanov [68] it became apparent that the “intermediate” cases require more complex log-gases models for their asymptotic descriptions, mixing both attracting and repelling charges, as well as allowing for charges “flowing” from one ensemble to another, see Sections 1.2 and 2 below.

In the context of the non-hermitian orthogonality, multiple orthogonal polynomials yield naturally the notion of the *critical vector-valued measures*, defined again as saddle points (and thus, solutions of some variational problems) for the total logarithmic energy.

In the same vein as the standard non-hermitian orthogonality can be characterized by a 2×2 matrix valued Riemann-Hilbert problem [28, 36], the multiple orthogonality can be associated with larger ($k \times k$, $k \geq 3$) matrix valued Riemann-Hilbert problems [37]. Some heuristics (that in some situations can be made rigorous using the WKB approximation arguments [61]) indicates that the asymptotics of multiple orthogonal polynomials (and the critical vector-valued measures) can be described

using higher-order algebraic equations. This is indeed the case, as it follows from some previous works and from the results of this paper.

Some instances of the critical vector-valued measures have appeared in literature, basically in the study of the limit zero distribution of some multiple orthogonal polynomials satisfying non-hermitian orthogonality conditions. The most studied cases are those presenting strong symmetry, so that the support of these measures can be easily derived. Other works, such as [6, 7], consider non-symmetric situations. Their basic approach is to start from a general form expected for a cubic equation that should govern the system, and then make a genus ansatz, hence reducing the free parameters of the equation to a minimum. Although a vector equilibrium is present, due to the explicit form of the deduced algebraic equation its detailed study is not needed and usually bypassed. Moreover, in such situations there is no external field involved, which is indeed a considerable simplification, even in the scalar case.

One of our main motivations is to understand the structure of the critical vector-valued measures with harmonic external fields. We follow the approach initiated in [53, 58], and show that the variational method, applied to critical measures, allows to deduce the corresponding algebraic equation (spectral curve) and the associated trajectories of a quadratic differential, now on a compact Riemann surface instead of the plane, whose analysis will be the central part of this paper.

1.2. Main results: general polynomial external fields. Assume we are given a vector of three non-negative measures $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$, compactly supported on the plane. For the interaction matrix A and vector of real-valued functions (“external fields”) $\vec{\phi} = (\phi_1, \phi_2, \phi_3)$

$$A = (a_{j,k}) = \begin{pmatrix} 1 & \frac{1}{2} & \frac{1}{2} \\ \frac{1}{2} & 1 & -\frac{1}{2} \\ \frac{1}{2} & -\frac{1}{2} & 1 \end{pmatrix}, \quad \phi_j = \operatorname{Re} \Phi_j, \quad j = 1, 2, 3, \quad (1.1)$$

where Φ_1, Φ_2, Φ_3 are polynomials¹, and denoting by

$$I(\mu, \nu) := \iint \log \frac{1}{|x - y|} d\mu(x) d\nu(y)$$

the logarithmic energy of two compactly supported measures μ and ν , we consider the total energy functional of the form [39]

$$E(\vec{\mu}) = E(\vec{\mu}, \vec{\phi}) = \sum_{j,k=1}^3 a_{j,k} I(\mu_j, \mu_k) + \sum_{j=1}^3 \int \phi_j d\mu_j. \quad (1.2)$$

Matrix A is symmetric, singular and positive-semidefinite. Its only non-zero eigenvalue (of order 2) is $3/2$, and the corresponding eigenspace is given by the vectors $(v_1, v_2, v_3)^T \in \mathbb{C}^3$ (in what follows, $(\cdot)^T$ stands for the transpose) satisfying

$$v_1 - v_2 - v_3 = 0.$$

¹Although our methodology is, generally speaking, valid also for general positive-semidefinite interaction matrices and when Φ_j are rational, we restrict ourselves here to the given matrix A and the polynomial external fields; these choices avoid some technicalities related to cumbersome notation and the order or coincidence of the poles.

As it will be seen below, the external fields in (1.1) must satisfy the compatibility condition of $\vec{\Phi}' = (\Phi'_1, \Phi'_2, \Phi'_3)^T$ being an eigenvector for the eigenvalue $3/2$, that is,

$$\Phi'_1(z) - \Phi'_2(z) = \Phi'_3(z), \quad z \in \mathbb{C}. \quad (1.3)$$

We complete the description of the electrostatic model by introducing the constraints on vector-valued measures under consideration: given a parameter $\alpha \in [0, 1]$, we require that $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$ satisfies

$$|\mu_1| + |\mu_2| = 1, \quad |\mu_1| + |\mu_3| = \alpha, \quad |\mu_2| - |\mu_3| = 1 - \alpha. \quad (1.4)$$

Here $|\mu|$ denotes the total mass (total variation) of the measure μ .

It should be pointed out that a particular case of this model (with all $\Phi_j \equiv 0$ and $\alpha = 1/2$) appears in the work of Aptekarev [3] and Rakhmanov [68] in their study of Hermite-Padé approximants [6, 8, 70]. Notice that A in (1.1) contains submatrices that are the interaction matrices both for the Angelesco and for the Nikishin cases.

Let us denote by mes_2 the plane Lebesgue measure on \mathbb{C} . We are interested in the existence of critical vector-valued measures within a family \mathcal{M}_α of measures $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$ satisfying (1.4) plus the following conditions:

- each μ_j , $j = 1, 2, 3$, is a non-negative Borel measure, supported on a compact set in \mathbb{C} of zero plane Lebesgue measure,

$$\text{mes}_2(\text{supp } \mu_j) = 0, \quad j = 1, 2, 3,$$

and such that the energy $E(\vec{\mu})$, defined in (1.2), is finite.

- the set

$$S_\alpha := \bigcup_{1 \leq j < k \leq 3} (\text{supp } \mu_j \cap \text{supp } \mu_k) \quad (1.5)$$

is finite,

$$\text{card } S_\alpha < \infty. \quad (1.6)$$

The *vector equilibrium problems* deal with the minimizers of the energy functional (1.2) over a family of measures $\vec{\mu}$ living on prescribed (and fixed) subsets of \mathbb{C} , see e.g. [13, 39, 41] and the references therein. As it follows from [13], the equilibrium measure (global minimizer of $E(\vec{\mu})$), if it exists in the class \mathcal{M}_α , is unique.

Here we consider a natural extension of the notion of vector equilibrium for the conducting plane: the *critical vector-valued measures*, defined as follows. For $t \in \mathbb{C}$ and $h \in C^2(\mathbb{C})$, denote by μ^t the pushforward measure of μ induced by the variation of the plane $z \mapsto h_t(z) = z + th(z)$, $z \in \mathbb{C}$.

Definition 1.1. We say that for $\alpha \in [0, 1]$, $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$ is an α -critical measure (or a saddle point of the energy $E(\cdot)$) if $\vec{\mu} \in \mathcal{M}_\alpha$, and

$$\lim_{t \rightarrow 0} \frac{E(\vec{\mu}^t) - E(\vec{\mu})}{t} = 0, \quad (1.7)$$

for every function $h \in C^2(\mathbb{C})$.

Sometimes, when the value of α is irrelevant or is clear from the context, we will simplify the terminology speaking simply about *critical measures*.

In order to formulate our first results we introduce some notation. Given a (scalar) non-negative Borel measure λ , compactly supported on the plane, we denote by

$$U^\lambda(z) := \int \log \frac{1}{|x - z|} d\lambda(x)$$

its logarithmic potential, which is harmonic in $\mathbb{C} \setminus \text{supp } \lambda$ and superharmonic in \mathbb{C} . Additionally, the principal value of the Cauchy transform

$$C^\lambda(z) := \lim_{\epsilon \rightarrow 0+} \int_{|z-x|>\epsilon} \frac{1}{x-z} d\lambda(x) \quad (1.8)$$

is analytic in $\mathbb{C} \setminus \text{supp } \lambda$, and

$$2 \frac{\partial U^\lambda}{\partial z}(z) = C^\lambda(z), \quad z \in \mathbb{C} \setminus \text{supp } \lambda, \quad (1.9)$$

with

$$\frac{\partial}{\partial z} = \frac{1}{2} \left(\frac{\partial}{\partial x} - i \frac{\partial}{\partial y} \right).$$

The interaction matrix A in (1.1) satisfies the following identities:

$$2A = B^T B = 3I_3 - bb^T, \quad B = \begin{pmatrix} 1 & 1 & 0 \\ -1 & 0 & -1 \\ 0 & -1 & 1 \end{pmatrix}, \quad b = \begin{pmatrix} 1 \\ -1 \\ -1 \end{pmatrix}, \quad (1.10)$$

where I_3 is the 3×3 identity matrix. For a vector of measures $\vec{\mu} = (\mu_1, \mu_2, \mu_3) \in \mathcal{M}_\alpha$, define the vector of functions $\vec{\xi} = (\xi_1, \xi_2, \xi_3)^T$ through the equality

$$\vec{\xi} = B \left(\frac{1}{3} \vec{\Phi}' + \vec{C} \right), \quad (1.11)$$

where $\vec{C} = (C^{\mu_1}, C^{\mu_2}, C^{\mu_3})^T$. Alternatively, the functions ξ_1, ξ_2, ξ_3 are given explicitly as

$$\begin{aligned} \xi_1(z) &= \frac{\Phi'_1(z)}{3} + \frac{\Phi'_2(z)}{3} + C^{\mu_1}(z) + C^{\mu_2}(z), \\ \xi_2(z) &= -\frac{\Phi'_1(z)}{3} - \frac{\Phi'_3(z)}{3} - C^{\mu_1}(z) - C^{\mu_3}(z), \\ \xi_3(z) &= -\frac{\Phi'_2(z)}{3} + \frac{\Phi'_3(z)}{3} - C^{\mu_2}(z) + C^{\mu_3}(z). \end{aligned} \quad (1.12)$$

It was proved in [58, Section 5] (see also [53]) that the Cauchy transform of a scalar critical measure is a solution of a quadratic equation mes_2 -a.e. on \mathbb{C} , which is a characterizing property of such measures. Our first result is a natural generalization of this property to the critical vector-valued measures:

Theorem 1.2. *Suppose that for $\alpha \in [0, 1]$, $\vec{\mu} = (\mu_1, \mu_2, \mu_3) \in \mathcal{M}_\alpha$ is an α -critical measure in the sense of Definition 1.1. Then there exist a polynomial R and a rational function D with poles of order at most 2 such that the functions ξ_1, ξ_2, ξ_3 in (1.11) satisfy mes_2 -a.e. on \mathbb{C} the algebraic equation*

$$\xi^3 - R(z)\xi + D(z) = 0. \quad (1.13)$$

The polynomial R is given by

$$R(z) = \frac{1}{9} \sum_{j,k=1}^3 a_{j,k} \Phi'_j(z) \Phi'_k(z) - \sum_{j=1}^3 \int \frac{\Phi'_j(x) - \Phi'_j(z)}{x-z} d\mu_j(x), \quad z \in \mathbb{C}. \quad (1.14)$$

Moreover, each of the measures μ_1, μ_2, μ_3 is supported on a finite union of analytic arcs, and they are absolutely continuous with respect to the arclength measure of their supports.

Each pole p of D , if exists, belongs to the support of at least two of the measures μ_1, μ_2, μ_3 , and for each such a measure μ_k its density satisfies

$$\frac{d\mu_k}{|dz|}(z) = |z - p|^{-\nu}(c + o(1)), \quad \text{as } z \rightarrow p \text{ along } \text{supp } \mu_k, \quad (1.15)$$

for some nonzero constant c (which may depend on p and k), where 3ν is the order of p as a pole of D .

We also prove a converse to Theorem 1.2.

Theorem 1.3. *Suppose that for $\alpha \in [0, 1]$, $\vec{\mu} = (\mu_1, \mu_2, \mu_3) \in \mathcal{M}_\alpha$ is such that the measures μ_1, μ_2, μ_3 are supported on a finite union of analytic arcs, and are absolutely continuous with respect to the arclength.*

If for some polynomials V_1, V_2, V_3 with

$$V_1'(z) + V_2'(z) + V_3'(z) = 0, \quad z \in \mathbb{C}, \quad (1.16)$$

the functions

$$\begin{aligned} \xi_1(z) &= V_1'(z) + C^{\mu_1}(z) + C^{\mu_2}(z), \quad z \in \mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2), \\ \xi_2(z) &= V_2'(z) - C^{\mu_1}(z) - C^{\mu_3}(z), \quad z \in \mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_3), \\ \xi_3(z) &= V_3'(z) - C^{\mu_2}(z) + C^{\mu_3}(z), \quad z \in \mathbb{C} \setminus (\text{supp } \mu_2 \cup \text{supp } \mu_3). \end{aligned} \quad (1.17)$$

satisfy an algebraic equation of the form (1.13) for a polynomial R and a rational function D , then $\vec{\mu}$ is a critical measure for the energy functional (1.1)–(1.2) and the external fields defined by

$$\Phi_1(z) := V_1(z) - V_2(z), \quad \Phi_2(z) := V_1(z) - V_3(z), \quad \Phi_3(z) := V_3(z) - V_2(z). \quad (1.18)$$

Notice that with the definition (1.18) and the compatibility condition (1.16) the equalities (1.17) take the form (1.12).

Few remarks are in order.

Theorems 1.2 and 1.3 show that there is a clear connection between critical measures and cubic equations of a specific form. The fact that the equation is cubic has to do with the rank of the interaction matrix A in (1.1), which is 2: from the proofs it is natural to expect that critical measures (with ≥ 3 components) for an energy functional with an interaction matrix of rank r should be related to algebraic equations of degree $r + 1$.

Furthermore, the compatibility condition (1.16) has to do with the structure of the interaction matrix (1.1), and finds its expression also in the absence of the ξ^2 term in the algebraic equation (1.13). Equation (1.16) and the structure of (1.17) should be seen as normalization conditions, ensuring that the underlying energy functional is of the precise form (1.2). Different choices for these identities would lead to a different interaction matrix in (1.2).

The notion of the critical vector-valued measures can be extended to accommodate also log-rational external fields and domains with non-trivial boundaries (“hard edges”, in the terminology of random matrix theory). The variations of the energy in these situations must allow for fixed points, see [58] for the scalar case. Fixed points and singularities of the external fields will appear among possible points of blow-up of the density of the critical measures, and also as possible poles of the coefficients R and D in (1.13).

Finally, a reader familiar with the notion of the scalar equilibrium and scalar critical measures might be puzzled by the fact that we allow for a blow-up of the

densities of such measures in the case of the polynomial external fields. Indeed, in the scalar case all critical measures in such circumstances have uniformly bounded densities on the whole complex plane [53]. This is no longer the case in the vector case, as the following two examples illustrate: they show that the components of the critical vector-valued measures can create “artificial hard edges”, obstructing each other and presenting unbounded densities when it was not expected, i.e. not created by the external field or by the boundaries of the conducting domain.

Example 1.4 (Repulsive or Angelesco interaction). Consider the algebraic equation

$$\xi^3 - \xi - \frac{1}{z} = 0.$$

It has simple branch points at $\pm 3\sqrt{3}/2$, and a double branch point at $z = 0$, and three solutions ξ_j specified by their behavior as $z \rightarrow \infty$:

$$\xi_1(z) = -\frac{1}{z} + \mathcal{O}\left(\frac{1}{z^2}\right), \quad \xi_2(z) = 1 + \frac{1}{2z} + \mathcal{O}\left(\frac{1}{z^2}\right), \quad \xi_3(z) = -1 + \frac{1}{2z} + \mathcal{O}\left(\frac{1}{z^2}\right),$$

with $\xi_1(z) + \xi_2(z) + \xi_3(z) = 0$. A direct calculation shows that there are two measures, μ_1 , living on $[-3\sqrt{3}/2, 0]$, and $\mu_2(x) = \mu_1(-x)$, such that $|\mu_1| = |\mu_2| = 1/2$, and

$$\begin{aligned} C^{\mu_1}(z) &= 1 - \xi_2(z), \quad z \in \mathbb{C} \setminus [-3\sqrt{3}/2, 0], \\ C^{\mu_2}(z) &= -1 - \xi_3(z), \quad z \in \mathbb{C} \setminus [0, 3\sqrt{3}/2]. \end{aligned}$$

Hence Theorem 1.3 shows that $\vec{\mu} = (\mu_1, \mu_2, 0) \in \mathcal{M}_{1/2}$ is a $1/2$ -critical vector-valued measure, according to Definition 1.1, corresponding to the interaction matrix (1.1) and the polynomial external fields given by

$$\Phi_1(z) = -z = -\Phi_2(z), \quad \Phi_3(z) = -2z.$$

Notice however that the densities of both measures μ_j , $j = 1, 2$, blow up at $z = 0$ as $|z|^{-1/3}$.

It is worth mentioning that the same pair of measures (μ_1, μ_2) solves the Angelesco-type vector *equilibrium* problem of minimizing (1.2) under the assumptions

$$\text{supp } \mu_1 \subset (-\infty, 0], \quad \text{supp } \mu_2 \subset [0, +\infty), \quad |\mu_1| = |\mu_2| = \frac{1}{2}, \quad \mu_3 = 0.$$

This equilibrium problem does exhibit a hard edge at the origin, and the proof is based on the same type of energy variation, but leaving $z = 0$ as a fixed point. It should be pointed out that a similar problem (but on finite intervals) was analyzed first by Kalyagin in [44].

Example 1.5 (Attractive or Nikishin interaction). Let us consider now the following cubic equation, which (up to a certain rescaling) have appeared already in [11]:

$$\xi^3 - \frac{1}{3}\xi - \frac{1}{z^2} + \frac{2}{27} = 0. \tag{1.19}$$

As in the Example 1.4, it has simple branch points at $\pm 3\sqrt{3}/2$, and a double branch point at $z = 0$.

Let $(81 - 12z^2)^{1/2}$ denote the branch of the square root in $\mathbb{C} \setminus [-3\sqrt{3}/2, 3\sqrt{3}/2]$ with positive boundary values on $(-3\sqrt{3}/2, 3\sqrt{3}/2)$ from the upper half plane, and

let

$$\Upsilon(z) := \left(-1 + \frac{3(81 - 12z^2)^{1/2}}{2z^2} \right)^{1/3}$$

also have positive boundary values on the same interval from the upper half plane. Then Υ is holomorphic in $\mathbb{C} \setminus [-3\sqrt{3}/2, 3\sqrt{3}/2]$, and

$$\xi_1(z) = \frac{1}{3} \left(e^{\pi i/3} \Upsilon(z) + \frac{1}{e^{\pi i/3} \Upsilon(z)} \right) = \frac{1}{3} - \frac{1}{z} + \mathcal{O}\left(\frac{1}{z^2}\right) \text{ as } z \rightarrow \infty$$

is holomorphic in $\mathbb{C} \setminus [0, 3\sqrt{3}/2]$. Moreover, Cardano's formula shows that ξ_1 is a solution of (1.19); the other two solutions are given by

$$\xi_2(z) = \xi_1(-z) \quad \text{and} \quad \xi_3(z) = -\xi_1(z) - \xi_2(z).$$

For the selected branch of Υ , the $\xi_{1,\pm}$ boundary values of ξ_1 on $(0, 3\sqrt{3}/2)$ from the upper and lower half plane, respectively, satisfy

$$w(x) := \frac{1}{2\pi i} (\xi_{1,+} - \xi_{1,-})(x) > 0, \quad x \in (0, 3\sqrt{3}/2).$$

Thus, $d\mu_2(x) := w(x)dx$ is a positive unit Borel measure on $[0, 3\sqrt{3}/2]$, and its Cauchy transform is

$$C^{\mu_2}(z) = \xi_1(z) - \frac{1}{3}, \quad z \in \mathbb{C} \setminus [0, 3\sqrt{3}/2].$$

Analogously, $\mu_3(z) := \mu_2(-z)$ is a positive unit Borel measure on $[-3\sqrt{3}/2, 0]$, and

$$C^{\mu_3}(z) = -\xi_2(z) + \frac{1}{3}, \quad z \in \mathbb{C} \setminus [-3\sqrt{3}/2, 0].$$

Theorem 1.3 shows that $\vec{\mu} = (0, \mu_2/2, \mu_3/2) \in \mathcal{M}_{1/2}$ is a $1/2$ -critical vector-valued measure, according to Definition 1.1, corresponding to the interaction matrix (1.1) and the polynomial external fields given by

$$\Phi_1(z) = 0, \quad \Phi_2(z) = z = -\Phi_3(z).$$

Notice that now the densities of both measures μ_j , $j = 2, 3$, blow up at $z = 0$ as $|z|^{-2/3}$.

The next example illustrates that in Theorem 1.3 we cannot discard the degenerate cases when the coefficient D in (1.13) is identically zero, or other situations when (1.13) becomes reducible over the field of rational functions:

Example 1.6. If μ is a scalar critical measure for the external field $\operatorname{Re} V$, where V is a polynomial, then there exists a polynomial Q for which [53, 58]

$$\left(C^\mu(z) + \frac{1}{2} V'(z) \right)^2 = Q(z), \quad \text{mes}_2 - \text{a.e. } z \in \mathbb{C}. \quad (1.20)$$

If we set

$$\mu_1 = \mu_3 = 0, \quad \mu_2 = \mu, \quad \Phi_2 = V, \quad \Phi_1 = -\Phi_3 = \frac{V}{2},$$

then

$$E(\vec{\mu}^t) = I(\mu^t, \mu^t) + \int \operatorname{Re} V d\mu^t, \quad (1.21)$$

which shows that the vector-valued measure $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$ is α -critical for $\alpha = 0$. The function ξ_2 in (1.12) becomes zero, whereas ξ_1, ξ_3 are analytic continuations of

each other across $\text{supp } \mu$. From the proof of Theorem 1.2 in Section 2 we conclude that the corresponding algebraic equation (3.9) is then reducible and has the form

$$\xi(\xi^2 - Q(z)) = 0, \quad (1.22)$$

where Q is the polynomial in (1.20).

Reciprocally, if the Cauchy transform of a probability measure μ satisfies an algebraic equation of the form (1.20) for some polynomials V and Q , then the construction just carried out combined with Theorem 1.3 gives us that μ is a scalar critical measure for the external field $\text{Re } V$.

In very much the same spirit, if we now set

$$\mu_2 = \mu_3 = \mu, \quad \mu_1 = 0, \quad \Phi_1 = 2\Phi_2 = 2\Phi_3 = V, \quad (1.23)$$

then the energy of $\vec{\mu}^t$ is also given by (1.21), showing again that $\vec{\mu}^t$ is α -critical, now for $\alpha = 1$. The respective functions ξ_1, ξ_2 are analytic continuation of each other across $\text{supp } \mu$, whereas $\xi_3 \equiv 0$ on \mathbb{C} , and the corresponding algebraic equation is reducible to the form (1.22).

The critical measure in (1.23) violates the assumption (1.6), but the conclusions of Theorem 1.2 still hold. However, condition (1.6) cannot be simply dropped, as the following example shows:

Example 1.7. Let Φ_1, Φ_2 be two arbitrary polynomials. For $\alpha = 1/2$ and the polynomial vector external field $\vec{\phi} = (\phi_1, \phi_2, \phi_1 - \phi_2)^T$, $\phi_j = \text{Re } \Phi_j$, satisfying (1.3), consider the vector-valued measure $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$, with

$$\mu_1 = \mu_2 = \frac{1}{2}\mu, \quad \mu_3 = 0.$$

The total energy of $\vec{\mu}$ is

$$E(\vec{\mu}) = \frac{3}{4} \left(I(\mu, \mu) + \frac{2}{3} \int (\phi_1 + \phi_2) d\mu \right).$$

Thus, if μ is a *scalar* probability critical measure in the external field

$$\phi = \frac{2}{3}(\phi_1 + \phi_2).$$

then (1.7) for $\vec{\mu}$ takes place. In other words, $\vec{\mu}$ is $1/2$ -critical in the external field $\vec{\phi}$. However, the statement of Theorem 1.2 does not hold: each of the corresponding functions ξ_1, ξ_2, ξ_3 in (1.12) satisfies an algebraic equation of degree 2, but if $\phi_1 \neq \phi_2$, these equations cannot be combined into a single algebraic equation of degree 3. To be more precise, calculations in the spirit of the proof of Theorem 1.2 in Section 2 below show that if Q is the polynomial in the right hand side of the equation (1.20) for μ (with $V' = (\Phi'_1 + \Phi'_2)/3$), then

$$\xi_1^2 = Q, \quad \xi_2 = -\frac{1}{2}(\xi_1 + \Phi'_1 - \Phi'_2), \quad \xi_3 = -\frac{1}{2}(\xi_1 - \Phi'_1 + \Phi'_2).$$

Clearly, the only assumption not satisfied by the measure $\vec{\mu}$ in this example is (1.6).

A consequence of Theorem 1.2 is that the set

$$\Xi_\alpha = (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3) \setminus S_\alpha \quad (1.24)$$

is a finite union of disjoint analytic arcs, each of them contained in exactly one of the supports $\text{supp } \mu_1, \text{supp } \mu_2, \text{supp } \mu_3$.

Any orientation of an arc $\Sigma \in \Xi_\alpha$ defines its left and right sides, as well as the left (or $+$) and right (or $-$) continuous boundary values of each ξ_j on Σ , that we denote by $\xi_{j\pm}$. In what follows, we use extensively (and without further warning) the notation mod 3, so that $\xi_0 = \xi_3$, $\xi_4 = \xi_1$ and so forth.

Since ξ_j 's are the three branches of the cubic equation (1.13), assumption (1.6) implies that for any curve $\Sigma \subset \Xi_\alpha$ there exists an index $j \in \{1, 2, 3\}$ such that

$$\xi_{j+}(z) = \xi_{j-}(z), \quad \xi_{(j+1)\pm}(z) = \xi_{(j-1)\mp}(z), \quad z \in \Sigma. \quad (1.25)$$

Critical measures are intimately connected with vector equilibrium problems and the vector symmetry (S -property), as it is evidenced by the next theorem.

Theorem 1.8. *Given $j \in \{1, 2, 3\}$, let Σ be an open analytic arc in $\text{supp } \mu_j \setminus S_\alpha$, not containing any branch point of (1.13). There exists a constant $l = l(\Sigma) \in \mathbb{R}$ for which both the Euler-Lagrange variational equation (equilibrium condition)*

$$\sum_{k=1}^3 a_{j,k} U^{\mu_k}(z) + \frac{\phi_j(z)}{2} = l, \quad (1.26)$$

and the S -property

$$\frac{\partial}{\partial n_+} \left(\sum_{k=1}^3 a_{j,k} U^{\mu_k}(z) + \frac{\phi_j(z)}{2} \right) = \frac{\partial}{\partial n_-} \left(\sum_{k=1}^3 a_{j,k} U^{\mu_k}(z) + \frac{\phi_j(z)}{2} \right)$$

hold true for $z \in \Sigma$, where n_\pm are the unit normal vectors to Σ , pointing in the opposite directions.

If, for $j = 1, 2, 3$, the constants $l(\Sigma) =: l_j$ above are independent of the connected component Σ of $\text{supp } \mu_j$, and additionally satisfy $l_1 - l_2 = l_3$, then the critical vector-valued measure is the equilibrium measure for the vector of contours $(\text{supp } \mu_1, \text{supp } \mu_2, \text{supp } \mu_3)$, see [13, Theorem 1.8]. Examples of critical but non-equilibrium measures in the scalar case are well-known (see e.g. [58]). Using the construction of Example 1.6, in particular (1.23), we can easily build vector critical non-equilibrium measures.

The S -property is the natural generalization of the scalar S -property, taking into account that in those arcs the potentials of the measures μ_{j-1} and μ_{j+1} are harmonic.

Critical vector-valued measures $\vec{\mu}$ are also closely connected with trajectories of a quadratic differential on the Riemann surface $\mathcal{R} = \mathcal{R}_1 \cup \mathcal{R}_2 \cup \mathcal{R}_3$ of the algebraic equation (1.13), with

$$\begin{aligned} \mathcal{R}_1 &= \overline{\mathbb{C}} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2), \\ \mathcal{R}_2 &= \overline{\mathbb{C}} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_3), \\ \mathcal{R}_3 &= \overline{\mathbb{C}} \setminus (\text{supp } \mu_2 \cup \text{supp } \mu_3). \end{aligned} \quad (1.27)$$

The sheets are connected across arcs in Ξ_α according to the following rule: if for $\Sigma \subset \Xi_\alpha$ the condition (1.25) holds, then we connect the sheets \mathcal{R}_{j+1} and \mathcal{R}_{j-1} through Σ crosswise, identifying the \pm -side of Σ on \mathcal{R}_{j+1} with its \mp -side on \mathcal{R}_{j-1} .

As usual, we regard the solutions ξ_1, ξ_2, ξ_3 in (1.12) as branches of the same meromorphic function $\xi : \mathcal{R} \rightarrow \mathbb{C}$, defined by (1.13), so that function ξ_j is the restriction of ξ to the sheet \mathcal{R}_j of the Riemann surface \mathcal{R} .

Theorem 1.9. *Let*

$$Q(z) = \begin{cases} \xi_2(z) - \xi_3(z), & \text{on } \mathcal{R}_1, \\ \xi_1(z) - \xi_3(z), & \text{on } \mathcal{R}_2, \\ \xi_1(z) - \xi_2(z), & \text{on } \mathcal{R}_3. \end{cases}$$

Then Q^2 extends to a meromorphic function on \mathcal{R} . Moreover,

$$\varpi = -Q^2(z)dz^2 \tag{1.28}$$

is a meromorphic quadratic differential on \mathcal{R} with possible poles only at the points at ∞ . Finally, for $j = 1, 2, 3$, each arc of $\text{supp } \mu_j$ is an arc of trajectory of ϖ .

As it was mentioned in Section 1.1, the connection between quadratic differentials and critical (or equilibrium) measures is not new. For scalar critical measures, the quadratic differential can be globally projected on the complex plane, so it is not necessary to consider it on a non-trivial Riemann surface. Even for vector-valued measures, this connection has been exploited before, see for instance [51, 54], but only locally on the plane. To our knowledge, Theorem 1.9 is the first time this connection is stated on the whole surface \mathcal{R} , which provides a powerful tool for construction of critical measures.

Remark 1.10. Assume that at $z = p$ we have (1.15) with $\nu \in \{1/3, 2/3\}$ (see Theorem 1.2). Then p is a triple branch point of \mathcal{R} , so that the local coordinate for \mathcal{R} at $z = p$ is $z = p + w^3$. Theorem 1.9 asserts that $\text{supp } \mu_j$ lives on trajectories of the quadratic differential ϖ defined in (1.28), which in this local coordinate takes the form $\varpi = w^{4-6\nu}(\text{const} + o(1))dw^2$, $\text{const} \neq 0$.

If $\nu = 1/3$ (or equivalently if the polynomial D in (1.13) has a simple pole at p), then by the local structure of the trajectories, their canonical projection from \mathcal{R} onto \mathbb{C} consists of two analytic arcs intersecting at p perpendicularly (a “cross”). As Example 1.4 shows, this configuration can correspond to the interaction of two repulsive measures, and p can potentially belong to the intersection of the three components μ_j .

However, for $\nu = 2/3$ (which means that D has a double pole at p) the canonical projection of the trajectories from \mathcal{R} onto \mathbb{C} at p consists of a single analytic arc through this point. It implies, due to (1.6), that in this case p belongs to the intersection of at most two (and then, exactly two) components of $\vec{\mu}$; moreover, it will be the endpoint of an arc in each of the supports. A quick analysis of the cubic branching shows that these two positive measures cannot interact repulsively (“Angelesco” interaction), and hence, must attract each other (“Nikishin” interaction). Thus, Example 1.5 represents (at least, locally) the most generic case of the $2/3$ -blowup of the density of a critical vector-valued measure.

Remark 1.11. We should point out that (1.6) plays a substantial role for Theorem 1.2, but the conclusion of Theorem 1.3 still holds true even if we drop (1.6) completely. In this situation, if Σ is an analytic arc contained in the support of at least two of the measures μ_1, μ_2, μ_3 , and $\mu_2 \neq \mu_3$ on Σ , then we can choose an orientation on Σ for which

$$\xi_{1+}(z) = \xi_{3-}(z), \quad \xi_{2+}(z) = \xi_{1-}(z), \quad \xi_{3+}(z) = \xi_{2-}(z), \quad z \in \Sigma. \tag{1.29}$$

The Euler-Lagrange equation (1.26) holds true for the three measures μ_1, μ_2, μ_3 (and possibly different constants l), and with respect to the orientation (1.29) the

S -property now becomes

$$\frac{\partial}{\partial n_+} \left(\sum_{k=1}^3 a_{jk} U^{\mu_k}(z) + \frac{\phi_j(z)}{2} \right) = (-1)^{j+1} \frac{\partial}{\partial n_-} \left(\sum_{k=1}^3 a_{j+1,k} U^{\mu_k}(z) + \frac{\phi_{j+1}(z)}{2} \right).$$

The conclusion of Theorem 1.9 is also valid in this situation (with appropriate gluing of the sheets of \mathcal{R} across the overlaps of the supports). However, we do not know whether there are positive measures with non-trivial overlapping supports for which the functions (1.12) are the three solutions to the same cubic equation as in Theorem 1.3.

1.3. Main results: the cubic case. In order to illustrate the results stated above we analyze in depth the interesting and highly non-trivial cubic case, corresponding to the energy functional (1.1)–(1.2), with

$$\Phi_1(z) = \Phi_2(z) = z^3, \quad \Phi_3(z) = 0, \quad z \in \mathbb{C}, \quad (1.30)$$

so that the external fields are

$$\phi(z) = \phi_1(z) = \phi_2(z) = \operatorname{Re} z^3, \quad \phi_3(z) = 0, \quad z \in \mathbb{C}; \quad (1.31)$$

in what follows, we consider only the values of the parameter $\alpha \in [0, 1/2]$; see the discussion in the Remark 1.15 and also in the introduction to Section 3.

We build a continuous one-parameter family of critical vector-valued measures $\vec{\mu}$ by choosing in (1.13) the coefficients

$$R(z) = 3z^4 - 3z - c, \quad D(z) = -2z^6 + 3z^3 + cz^2 - 3\alpha(1 - \alpha), \quad (1.32)$$

where $c = c(\alpha)$ is the real parameter given by

$$c = - \left(\frac{243}{64} (1 - 4\alpha(1 - \alpha))^2 \right)^{\frac{1}{3}}. \quad (1.33)$$

Theorem 1.12. *For $\alpha \in [0, 1/2]$, there exists a one-parameter family $\vec{\mu} = \vec{\mu}_\alpha \in \mathcal{M}_\alpha$ of critical measures for the potentials (1.31) for which the corresponding algebraic equation (1.13) has coefficients (1.32).*

Moreover, for any such a critical measure, if the associated algebraic equation (1.13) defines a Riemann surface of genus 0, then its coefficients are given by (1.32) for some choice of the cubic root in (1.33).

Remark 1.13. As it follows from the proof of this theorem in Section 3.1 (see also Remark 3.2), there are other choices for the coefficient c that yield critical vector-valued measures, but for which the associated Riemann surface is of genus 1. We do not consider this case here.

Theorem 3.7 below (see also Remark 3.8) shows that the measures $\vec{\mu}_\alpha$ from Theorem 1.12 are in fact equilibrium measures on an appropriate family of contours in \mathbb{C} .

For the critical measure $\vec{\mu}_\alpha$ given by Theorem 1.12, the algebraic function ξ defined by (1.13) has two real branch points $a_1 < b_1$ and two nonreal branch points $b_2, a_2 = \bar{b}_2$ ($\operatorname{Im} b_2 > 0$). The support of the components μ_j of the critical measure $\vec{\mu}_\alpha = (\mu_1, \mu_2, \mu_3)$ can be easily described; we show that it exhibits a phase transition for a certain value of α :

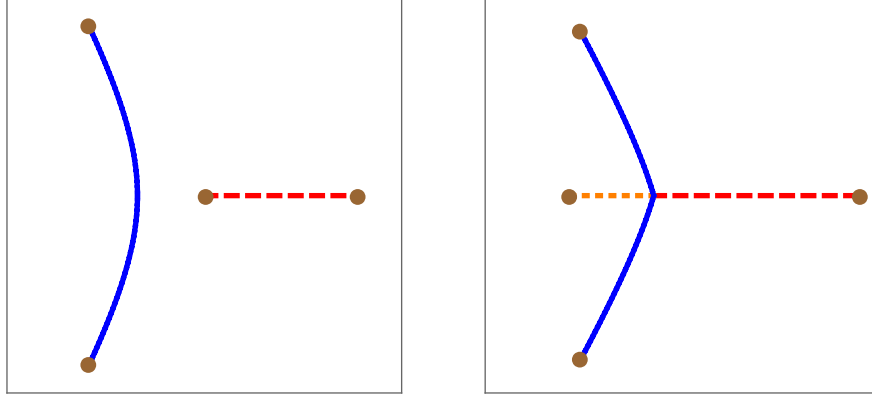


FIGURE 1. For $\tau \approx 0.126 < \tau_c$ (left panel) and for $\tau \approx 0.227 > \tau_c$ (right panel), numerical evaluation of the supports of the measures μ_1 (long dashed line), μ_2 (continuous line) and μ_3 (short dashed line - only on the right panel).

Theorem 1.14. *There exists a critical value $\alpha_c \in (0, 1/2)$, determined by*

$$\operatorname{Im} \int_{b_2}^{a_1} (\xi_1(s) - \xi_3(s)) ds = 0, \quad (1.34)$$

where ξ_1, ξ_3 are the functions in (1.12), such that

- If $\alpha < \alpha_c$ then $\mu_3 = 0$, $\operatorname{supp} \mu_1 = [a_1, b_1]$, and $\operatorname{supp} \mu_2$ is an analytic arc, disjoint from $\operatorname{supp} \mu_1$ and connecting the branch points a_2, b_2 .
- If $\alpha_c < \alpha < 1/2$ then none of the measures is zero, and there exists a value $a_* \in (a_1, b_1)$, determined by the condition

$$\operatorname{Im} \int_{b_2}^{a_*} (\xi_1(s) - \xi_3(s)) ds = 0, \quad (1.35)$$

for which $\operatorname{supp} \mu_1 = [a_*, b_1]$, $\operatorname{supp} \mu_3 = [a_1, a_*]$ and $\operatorname{supp} \mu_* = \gamma_1 \cup \gamma_2$, where γ_1 is an analytic arc on the upper half plane connecting b_2 to a_* , and γ_2 is its complex conjugate arc.

We refer to Figure 1 for an actual numerical evaluation of the supports of the measures μ_1, μ_2, μ_3 in Theorem 1.14.

Remark 1.15. The restriction $\alpha < 1/2$ arises somewhat naturally in our analysis. For $\alpha = 1/2$, the algebraic equation (1.13) with coefficients (1.32)–(1.33) still defines a vector-valued measure $\vec{\mu}_{1/2}$ in the sense of Theorem 1.3, but the support of its third component is unbounded; hence, strictly speaking the measure $\vec{\mu}_{1/2}$ does not belong to $\mathcal{M}_{1/2}$.

In rough terms, Theorems 1.2 and 1.3 induce a map

$$\{\text{critical vector-valued measures}\} \mapsto \{\text{algebraic equations}\}$$

and Theorems 1.12 and 1.14 characterize this map completely for a subclass of algebraic equations and the cubic potential (1.30).

The existence of the critical measure $\vec{\mu}$ assured by Theorem 1.12 and the topological description of its support given by Theorem 1.14 follow after a careful analysis of the associated quadratic differential ϖ in (1.28). We describe the critical graph of ϖ for $\alpha = 0$ and indicate its subsets corresponding to the supports $\text{supp } \mu_1, \text{supp } \mu_2, \text{supp } \mu_3$. In other words, we lift the supports to the associated Riemann surface, embedding them into the critical graph of ϖ . We then deform the parameter α in the interval $(0, 1/2)$, observing the dynamics of the critical graph and keeping track of the supports of the measures.

During this deformation procedure, the critical graph of ϖ displays several transitions, described in Sections 4.5.2–4.5.7, but only one of such transitions, governed by equation (1.34), has direct impact on the topology of $\text{supp } \mu_1, \text{supp } \mu_2, \text{supp } \mu_3$.

It should be stressed that there are not many tools for such an analysis (we have summarized the basic facts about quadratic differentials in the Appendix B), and that it can be extremely complicated, even on the Riemann sphere. The situation becomes even more involved when we consider trajectories on a more general compact Riemann surface. We believe that the methodology we have developed is of independent interest, and can be applied in other situations as well.

1.4. Critical measures and max-min problems. The critical measure given by Theorem 1.12 is crucial in the study of multiple orthogonal polynomial $P_{n,m}(z) = z^N + \dots$ defined through

$$\begin{aligned} \int_{\Sigma_1} z^j P_{n,m}(z) e^{-Nz^3} dz &= 0, \quad j = 0, \dots, n-1, \\ \int_{\Sigma_2} z^j P_{n,m}(z) e^{-Nz^3} dz &= 0, \quad j = 0, \dots, m-1, \end{aligned} \quad (1.36)$$

where $N = n + m$ and Σ_1, Σ_2 are contours extending to ∞ with angles $-\frac{2\pi}{3}, 0$ and $-\frac{2\pi}{3}, \frac{2\pi}{3}$, respectively. It turns out that the asymptotic behavior of $P_{n,m}$ as $m, n \rightarrow \infty$ in such a way that $n/N \rightarrow \alpha \in (0, 1/2)$, is governed by this critical vector-valued measures $\vec{\mu} \in \mathcal{M}_\alpha$: under this regime, the sequence of zero counting measures associated with $(P_{n,m})$ converges to the sum of the first two components of $\vec{\mu}$. The two extremal cases, $\alpha = 0$ and $\alpha = 1/2$, have been proved in [27] and [35], respectively, and the general case $\alpha \in (0, 1/2)$ will be addressed in a follow-up publication [59].

As we already mentioned, zeros of orthogonal polynomials on the real line are characterized as solutions to logarithmic equilibrium problems. The measure $\vec{\mu}$ given by Theorem 1.12, and hence the zeros of the multiple orthogonal polynomials (1.36), also admits a vector equilibrium problem characterization, as explained next.

Let \mathcal{T} be the class of contours Σ extending to ∞ with angles $-\frac{2\pi}{3}, \frac{2\pi}{3}$ and intersecting the real axis exactly once, say at $a_* = a_*(\Sigma)$. Denote by $\mathcal{M}_\alpha(\Sigma)$ the subset of measures $\vec{\nu} = (\nu_1, \nu_2, \nu_3) \in \mathcal{M}_\alpha$ that further satisfy

$$\text{supp } \nu_1 \subset [a_*, +\infty), \quad \text{supp } \nu_2 \subset \Sigma, \quad \text{supp } \nu_3 \subset (-\infty, a_*].$$

The α -equilibrium measure $\vec{\nu}_\alpha = \vec{\nu}_{\alpha, \Sigma}$ of $\Sigma \in \mathcal{T}$ (for the energy functional $E(\cdot)$ given in (1.2) with (1.30)–(1.31)) is the unique measure $\vec{\nu}_\alpha \in \mathcal{M}_\alpha(\Sigma)$ for which

$$E(\vec{\nu}_\alpha) = \inf_{\vec{\nu} \in \mathcal{M}_\alpha(\Sigma)} E(\nu) =: E_\alpha(\Sigma).$$

Theorem 1.16. *For $\alpha \in (0, 1/2)$, there exists a contour $\Gamma \in \mathcal{T}$ for which its α -equilibrium measure is the α -critical measure $\vec{\mu}$ given by Theorem 1.12.*

We conjecture that the contour Γ is characterized by the max-min property

$$E_\alpha(\Gamma) = \sup_{\Sigma \in \mathcal{T}} E_\alpha(\Sigma).$$

This characterization would be analogous to the max-min property for non-hermitian orthogonal polynomials [53, 69, 74, 75]. Furthermore, an adaptation of the methods in [53, 58] shows that the α -equilibrium measure of any contour solving this max-min problem is in fact an α -critical measure.

1.5. Structure of the rest of the paper. Theorems 1.2, 1.8, 1.9, and a simplified version of Theorem 1.3 are proven in Section 2 (the complete proof of Theorem 1.3 is rather technical and thus deferred to Appendix A). In Section 3.1 we derive the coefficients (1.32)–(1.33) and prove the second part of Theorem 1.12. Assuming the existence of certain short trajectories (fact proved later, in Section 4), we derive the first part of Theorem 1.12 and Theorem 3.7 in Section 3.2, which yields Theorem 1.16. Section 4 is devoted to the analysis of the global structure of the quadratic differential (1.28) for the cubic potentials (1.30), which also provides a proof of Theorem 1.14 (or equivalently, of Theorem 4.4, where the structure of the support of the critical vector-valued measure is rephrased in the terminology of quadratic differentials). Finally, two more appendices contain the minimal background on quadratic differentials (Appendix B) and the description of the numerical experiments and procedures used in this paper (Appendix C).

2. CRITICAL VECTOR-VALUED MEASURES

Critical measures were defined in the previous section in terms of the vanishing of the total energy (1.7). Here we use an alternative characterization, more convenient for calculations. When $\vec{\mu} = \mu$ is a scalar measure, the following Proposition 2.1 coincides with [58, Lemma 3.1]. The proof extends trivially to vectors of measures $\vec{\mu}$ and polynomial external fields as considered here. We skip the details.

Proposition 2.1. *For any vector of measures $\vec{\mu} \in \mathcal{M}_\alpha$, it is valid*

$$E(\vec{\mu}^t) = E(\vec{\mu}) + t \operatorname{Re} \mathcal{D}_h(\vec{\mu}) + \mathcal{O}(t^2), \quad t \rightarrow 0,$$

where

$$\mathcal{D}_h(\vec{\mu}) = - \sum_{j,k=1}^3 a_{j,k} \iint \frac{h(x) - h(y)}{x - y} d\mu_j(x) d\mu_k(y) + \sum_{j=1}^3 \int \Phi'_j(x) h(x) d\mu_j(x). \quad (2.1)$$

Although the definition of critical measures only requires to evaluate $\mathcal{D}_h(\vec{\mu})$ for test functions $h \in C^2(\mathbb{C})$, the quantity $\mathcal{D}_h(\vec{\mu})$ as defined above will be used for any function defined on $\operatorname{supp} \mu_1 \cup \operatorname{supp} \mu_2 \cup \operatorname{supp} \mu_3$, as long as the integrals on the right-hand side of (2.1) exist.

In particular, considering h and $-ih$, we get from Proposition 2.1

Corollary 2.2. *A measure $\vec{\mu} \in \mathcal{M}_\alpha$ is critical if, and only if,*

$$\mathcal{D}_h(\vec{\mu}) = 0,$$

for every function $h \in C^2(\mathbb{C})$.

The following lemma is inspired by [58], where a similar situation (but involving only one measure) was analyzed, and borrowed in the following form from [53, Lemma 3.5], where a simplified proof is given. The proof from [53] is easily extended to our situation by simply choosing $g \equiv 1$ therein and mimicking the arguments.

Lemma 2.3. *Suppose $\vec{\mu}$ is a critical vector-valued measure. If $z \in \mathbb{C}$ satisfies*

$$\int \frac{d\mu_j(x)}{|x-z|} < \infty, \quad j = 1, 2, 3, \quad (2.2)$$

then $\mathcal{D}_{h_z}(\vec{\mu}) = 0$ for the meromorphic function

$$h_z(x) = \frac{1}{x-z}, \quad x \in \mathbb{C}. \quad (2.3)$$

The next result is straightforward, and we skip its proof:

Lemma 2.4. *Let f be holomorphic in a punctured neighborhood of $p \in \mathbb{C}$. If $|f(z)|^r$ is mes_2 -locally integrable at p for some $r > 0$, then p is either a removable singularity or a pole of f of order $< 2/r$.*

For any finite Borel measure μ , it follows from Tonelli's theorem that the function

$$z \mapsto \int \frac{d\mu(x)}{|x-z|}$$

is locally mes_2 -integrable and in particular finite mes_2 -a.e.; that is, (2.2) is satisfied for mes_2 -a.e. $z \in \mathbb{C}$. Hence the Cauchy transform C^μ of the measure μ is finite mes_2 -a.e. In particular, the functions ξ_1, ξ_2, ξ_3 in (1.12) are finite mes_2 -a.e. and locally integrable with respect to mes_2 , so that

$$\sum_{j=1}^3 \xi_j^2(z) \quad (2.4)$$

is well-defined mes_2 -a.e. and is analytic on open subsets of $\mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3)$.

Recall the quantities R , $\mathcal{D}_h(\vec{\mu})$ and h_z defined in (1.14), (2.1) and (2.3), respectively. The next lemma is crucial for what comes later:

Lemma 2.5. *If $\vec{\mu} \in \mathcal{M}_\alpha$, then the functions ξ_j defined in (1.12) satisfy the identity*

$$\frac{1}{2} (\xi_1^2 + \xi_2^2 + \xi_3^2)(z) = R(z) + \mathcal{D}_{h_z}(\vec{\mu}) \quad \text{mes}_2\text{-a.e.} \quad (2.5)$$

Moreover, the statements

- (i) $\sum_{j=1}^3 \xi_j^2(z)$ is a polynomial,
- (ii) $\sum_{j=1}^3 \xi_j^2(z) = 2R(z)$ mes_2 -a.e.
- (iii) $\mathcal{D}_{h_z}(\vec{\mu}) = 0$ for every z satisfying (2.2),

are equivalent.

Proof. For $h = h_z$,

$$\begin{aligned} \iint \frac{h_z(x) - h_z(y)}{x-y} d\mu_j(x) d\mu_k(y) &= - \iint \frac{1}{(x-z)(y-z)} d\mu_j(x) d\mu_k(y) \\ &= -C^{\mu_j}(z) C^{\mu_k}(z), \end{aligned} \quad (2.6)$$

and

$$\int h_z(x) \Phi'_j(x) d\mu_j(x) = \Phi'_j(z) C^{\mu_j}(z) - Q_j(z), \quad (2.7)$$

where

$$Q_j(z) = - \int \frac{\Phi'_j(x) - \Phi'_j(z)}{x-z} d\mu_j(x). \quad (2.8)$$

Notice that the integrand in the right-hand side of (2.8) is a polynomial both in x and in z , and thus (since μ_j is finite and compactly supported) it is also $d\mu_j(x)$ -integrable for all z . Hence, Q_j is a bona fide polynomial. Additionally, (2.6) and (2.7) also show that $D_{h_z}(\vec{\mu})$ is finite whenever (2.2) holds true.

Using (2.6) and (2.7) in (2.1), we get that

$$\mathcal{D}_{h_z}(\vec{\mu}) = \sum_{j,k=1}^3 a_{j,k} C^{\mu_j}(z) C^{\mu_k}(z) + \sum_{j=1}^3 \Phi'_j(z) C^{\mu_j}(z) - \sum_{j=1}^3 Q_j(z),$$

or equivalently,

$$\mathcal{D}_{h_z}(\vec{\mu}) = \vec{C}(z)^T A \vec{C}(z) + (\vec{\Phi}'(z))^T \vec{C}(z) - \sum_{j=1}^3 Q_j(z), \quad (2.9)$$

with the previously used notation $\vec{C} = (C^{\mu_1}, C^{\mu_2}, C^{\mu_3})^T$ and $\vec{\Phi}' = (\Phi'_1, \Phi'_2, \Phi'_3)^T$.

On the other hand, by (1.11) (and taking also into account both identities from (1.10)) we get

$$\begin{aligned} \frac{1}{2} (\xi_1^2 + \xi_2^2 + \xi_3^2) &= \frac{1}{2} \xi^T \xi = \left(\frac{1}{3} \vec{\Phi}' + \vec{C} \right)^T A \left(\frac{1}{3} \vec{\Phi}' + \vec{C} \right) \\ &= \frac{1}{9} (\vec{\Phi}')^T A \vec{\Phi}' + \vec{C}^T A \vec{C} + \frac{2}{3} (\vec{\Phi}')^T A \vec{C} \\ &= \frac{1}{9} (\vec{\Phi}')^T A \vec{\Phi}' + \vec{C}^T A \vec{C} + \frac{1}{3} (\vec{\Phi}')^T (3I_3 - bb^T) \vec{C} \\ &= \left\{ \vec{C}^T A \vec{C} + (\vec{\Phi}')^T \vec{C} \right\} + \frac{1}{9} (\vec{\Phi}')^T A \vec{\Phi}' - \frac{1}{3} (b^T \vec{\Phi}')^T (b^T \vec{C}). \end{aligned}$$

Notice that the compatibility condition (1.3) means precisely that $b^T \vec{\Phi}' = \vec{0}$, so that

$$\frac{1}{2} (\xi_1^2 + \xi_2^2 + \xi_3^2) = \frac{1}{9} (\vec{\Phi}')^T A \vec{\Phi}' + \left\{ \vec{C}^T A \vec{C} + (\vec{\Phi}')^T \vec{C} \right\},$$

and using (2.9) we can simplify the expression between brackets above to conclude that

$$\frac{1}{2} (\xi_1^2 + \xi_2^2 + \xi_3^2) = \mathcal{D}_{h_z}(\vec{\mu}) + \frac{1}{9} (\vec{\Phi}')^T A \vec{\Phi}' + \sum_{j=1}^3 Q_j = \mathcal{D}_{h_z}(\vec{\mu}) + R,$$

where R was defined in (1.14). This identity is valid for values of z satisfying (2.2), which proves (2.5), as well as the equivalence of (ii) and (iii).

Clearly, (ii) implies (i). On the other hand, by (2.6),

$$\lim_{z \rightarrow \infty} \iint \frac{h_z(x) - h_z(y)}{x - y} d\mu_j(x) d\mu_k(y) = 0,$$

and since function $h_z(\cdot)$ converges uniformly to 0 as $z \rightarrow \infty$, by the Dominated Convergence Theorem we conclude that

$$\lim_{z \rightarrow \infty} \int h_z(x) \Phi'_j(x) d\mu_j(x) = 0.$$

Thus,

$$\lim_{z \rightarrow \infty} \mathcal{D}_{h_z}(\vec{\mu}) = 0. \quad (2.10)$$

Measures μ_1, μ_2 and μ_3 are compactly supported, so that (2.5) holds for all z sufficiently large, and with (2.10),

$$\lim_{z \rightarrow \infty} (\xi_1^2 + \xi_2^2 + \xi_3^2 - 2R)(z) = 0.$$

This establishes that (i) implies (ii). \square

Proof of Theorem 1.2. A combination of Corollary 2.2 and Lemma 2.5 shows that if $\vec{\mu}$ is critical, then the polynomial R in (1.14) is alternatively expressed as

$$R(z) = \frac{1}{2}(\xi_1(z)^2 + \xi_2(z)^2 + \xi_3(z)^2).$$

From the expressions (1.12) it follows that

$$\xi_1 + \xi_2 + \xi_3 = 0, \quad (2.11)$$

and as a consequence

$$R = \frac{1}{2}(\xi_1^2 + \xi_2^2 + \xi_3^2) = \xi_1^2 + \xi_1\xi_2 + \xi_2^2 = \xi_2^2 + \xi_2\xi_3 + \xi_3^2 = \xi_3^2 + \xi_3\xi_1 + \xi_1^2. \quad (2.12)$$

Multiplying the identities in (2.12) by appropriate factors $(\xi_j - \xi_k)$, it follows that

$$\xi_1(z)^3 - \xi_1(z)R(z) = \xi_2(z)^3 - \xi_2(z)R(z) = \xi_3(z)^3 - \xi_3(z)R(z). \quad (2.13)$$

Define

$$D(z) = -\xi_j(z)^3 + \xi_j(z)R(z). \quad (2.14)$$

Due to (2.13), D does not depend on the choice of $j \in \{1, 2, 3\}$. When we choose $j = 1$, we see that D is analytic on $\mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2)$. Similarly, when we choose $j = 2, 3$, we see that D should be analytic on $\mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_3)$ and $\mathbb{C} \setminus (\text{supp } \mu_2 \cup \text{supp } \mu_3)$, respectively. Consequently, D can only have singularities at S_α , defined in (1.5), so from (1.6) these singularities are all isolated and in a finite number.

Moreover,

$$\begin{aligned} |D(z)|^{1/3} &\leq (|\xi_j(z)|^3 + |\xi_j(z)||R(z)|)^{1/3} \\ &\leq |\xi_j(z)| + |\xi_j(z)|^{1/3}|R(z)|^{1/3}, \end{aligned}$$

where for the last inequality we used $(x + y)^{1/3} \leq x^{1/3} + y^{1/3}$, which is valid for any non negative numbers x, y .

The function ξ_j is locally mes_2 -integrable (see the comments after Lemma 2.4). Hence the function D satisfies the conditions of Lemma 2.4 for $r = 1/3$ and any choice of $p \in S_\alpha$, so the points in S_α are not essential singularities of D . Moreover, from the behavior of the ξ_j 's and R when $z \rightarrow \infty$, it follows that D has polynomial growth at ∞ . We conclude that D is rational (although it could be identically zero).

In summary, we have shown that under the assumptions of Theorem 1.2, there exist a polynomial R , defined in (1.14), and a rational function D with possible poles in S_α , such that the functions ξ_1, ξ_2, ξ_3 satisfy

$$\xi_j(z)^3 - R(z)\xi_j(z) + D(z) = 0, \quad \text{mes}_2\text{-a.e.}$$

The fact that all the components μ_j are supported on a finite union of analytic arcs, and that they are absolutely continuous with respect to the arclength measure of their supports is then a direct consequence of [19, Theorem 2].

Let p be a pole of D . The functions ξ_1, ξ_2, ξ_3 satisfy the cubic equation (1.13), so their local behavior near p is of the form

$$\xi_j(z) = (z - p)^{-\nu_j} (\kappa_j + o(1)), \quad \text{as } z \rightarrow p, \quad j = 1, 2, 3, \quad (2.15)$$

where $\nu_1, \nu_2, \nu_3 \in \mathbb{Q}$ and $\kappa_1, \kappa_2, \kappa_3$ are nonzero constants. Since p is a pole of D and R is a polynomial, from (2.14) we learn the ν_j 's are all equal, say to ν , and 3ν should be equal to the order of p as a pole of D . Moreover, we further get that none of the functions ξ_1, ξ_2 and ξ_3 is analytic near p , so p has to belong to the support of at least two of the measures μ_1, μ_2 and μ_3 . On the other hand, since ξ_j is a linear combination of polynomials and Cauchy transforms of finite measures, we must also have $\nu < 1$, so we conclude

$$\nu \in \left\{ \frac{1}{3}, \frac{2}{3} \right\}.$$

The behavior (1.15) then follows from Plemelj's formula and (2.15), concluding the proof. \square

Remark 2.6. The fact that (scalar) finite measures, whose Cauchy transform satisfies a quadratic equation mes₂-q.e., live on a finite set of analytic curves was established in [58] without any additional constraint on the measure (the proof was inspired by [14]). This was extended to an arbitrary algebraic equation in [19], but with a crucial assumption on the measure that its support has a zero plane Lebesgue measure (this assumption is embedded in the definition of the class \mathcal{M}_α in Section 1.2). Nevertheless, we believe that the assertion about the support is still true even if we drop this a priori restriction on $\vec{\mu}$.

Our next goal is to prove Theorem 1.3. In other words, for a vector of measures $\vec{\mu} \in \mathcal{M}_\alpha$ we assume that the components of $\vec{\mu}$ are supported on a finite union of analytic arcs, and the associated functions ξ_j in (1.17) satisfy an algebraic equation of the form (1.13) for some polynomial R and some rational function D . As we already pointed out after the statement of Theorem 1.3, these assumptions allow us to write ξ as in (1.12) for the vector of external fields $\vec{\Phi} = (\Phi_1, \Phi_2, \Phi_3)$ given by (1.18). In particular, it follows that ξ_j 's satisfy (2.11) and consequently

$$R = \xi_1 \xi_2 + \xi_1 \xi_3 + \xi_2 \xi_3 = \frac{1}{2}(\xi_1^2 + \xi_2^2 + \xi_3^2).$$

From Lemma 2.5 we get immediately the following result:

Corollary 2.7. *Under the assumptions of Theorem 1.3, the vector of measures $\vec{\mu}$ satisfies*

$$\mathcal{D}_{h_z}(\vec{\mu}) = 0$$

for $z \in \mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3)$, where h_z is as in (2.3).

Hence, the proof of Theorem 1.3 is reduced basically to extending the conclusion of Corollary 2.7 to the whole class of functions $h \in C^2(\mathbb{C})$. Although this is relatively straightforward for h analytic in a neighborhood of the union of $\text{supp } \mu_j$, it takes more effort beyond the analyticity due to the integrand $(h(x) - h(y))/(x - y)$ in (2.1). In this case we need to appeal to a “refined” version of the celebrated Mergelyan's theorem that would allow for an approximation of a given sufficiently smooth function by a sequence of analytic ones, and in such a way that the sequence of derivatives is uniformly bounded. The complete proof we were able to find is

rather technical, so for the sake of readability we postpone it to Appendix A, and present here the arguments valid for the case when h vanishes in a neighborhood of the set of double poles of D (or when this set is empty). We formulate it as an independent proposition:

Proposition 2.8. *Let $\vec{\mu} \in \mathcal{M}_\alpha$ be a vector of measures satisfying the conditions of Theorem 1.3. Suppose that $h \in C^2(\mathbb{C})$ satisfies*

$$\text{supp } h \cap \widehat{S}_\alpha = \emptyset,$$

where \widehat{S}_α is the set of double poles of D . Then $\mathcal{D}_h(\vec{\mu}) = 0$.

Proof. Recall that the union of the support of the components of $\vec{\mu}$ is a finite union of analytic arcs, whose points of intersection are denoted by S_α (see (1.5)), and the poles of the coefficient D in (1.13) belong to S_α , as it follows from the arguments in (2.15) *et seq.* Furthermore, if $p \in S_\alpha \setminus \widehat{S}_\alpha$, then

$$\frac{d\mu_j}{ds}(z) = \mathcal{O}(|z - p|^{-\nu}), \quad \text{as } z \rightarrow p, \quad j = 1, 2, 3, \quad (2.16)$$

for some $\nu \leq 1/3$ (possibly $\nu < 0$).

Given h as above, let us define

$$\vec{g}(z) = (g_1(z), g_2(z), g_3(z))^T \quad \text{and} \quad \vec{H}(z) = (H_1(z), H_2(z), H_3(z))^T,$$

with

$$g_j(z) = \int \frac{h(x)}{x - z} d\mu_j(x), \quad H_j(z) = \int \frac{h(x) - h(z)}{x - z} d\mu_j(x), \quad j = 1, 2, 3,$$

where as usual we understand $g_j(z)$ in terms of its principal value (as in (1.8)), for $z \in \text{supp } \mu_j$. With this convention and recalling (1.24), on any open analytic arc Σ from Ξ_α , not containing any branch point of (1.13), by the Sokhotsky–Plemelj’s formula we have

$$g_{j\pm}(z) = \pm \pi i h(z) \mu'_j(z) + g_j(z), \quad z \in \Sigma, \quad (2.17)$$

where $\mu'_j = d\mu_j/ds$ is the Radon-Nikodym derivative of μ_j with respect to the line element on Σ , so that if $\Sigma \cap \text{supp } \mu_j = \emptyset$, we have $\mu'_j = 0$.

Let Σ be an open analytic arc from $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$, not containing any branch point of (1.13). Observe that the matrix B in (1.10) satisfies the identity $B^T = M_1 - M_2$, where

$$M_1 = \begin{pmatrix} 1 & 0 & 0 \\ 0 & 0 & -1 \\ 0 & -1 & 0 \end{pmatrix}, \quad M_2 = \begin{pmatrix} 0 & 1 & 0 \\ -1 & 0 & 0 \\ 0 & 0 & -1 \end{pmatrix}.$$

Since \vec{H} is continuous across Σ , by (1.10),

$$2A\vec{H}(z) = (M_1 - M_2) B \vec{H}(z) = M_1 B \vec{H}_-(z) - M_2 B \vec{H}_+(z), \quad z \in \Sigma. \quad (2.18)$$

Furthermore, the identity

$$\vec{H}(z) = \vec{g}(z) - h(z) \vec{C}(z)$$

holds for all $z \in \mathbb{C} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3)$, and we obtain that

$$2A\vec{H}(z) = (M_1 B \vec{g}_- - M_2 B \vec{g}_+)(z) - h(z) (M_1 B \vec{C}_- - M_2 B \vec{C}_+)(z). \quad (2.19)$$

Recall that Equation (1.16) allows us to rewrite (1.17) as (1.12) with the external fields given in (1.18). From (1.11),

$$B\vec{C} = \vec{\xi} - \frac{1}{3}B\vec{\Phi}',$$

so that

$$\begin{aligned} M_1 B\vec{C}_-(z) - M_2 B\vec{C}_+(z) &= \left(M_1(\vec{\xi}_- - \frac{1}{3}B\vec{\Phi}') - M_2(\vec{\xi}_+ - \frac{1}{3}B\vec{\Phi}') \right)(z) \\ &= M_1\vec{\xi}_-(z) - M_2\vec{\xi}_+(z) - \frac{2}{3}A\vec{\Phi}'(z). \end{aligned} \quad (2.20)$$

Because the functions ξ_j 's satisfy a cubic equation, the boundary conditions (1.25) are valid and imply the equality

$$X M_1 \vec{\xi}_-(z) = X M_2 \vec{\xi}_+(z), \quad z \in \Sigma, \quad (2.21)$$

where

$$X = X(z) = \begin{pmatrix} \chi_1(z) & 0 & 0 \\ 0 & \chi_2(z) & 0 \\ 0 & 0 & \chi_3(z) \end{pmatrix},$$

and χ_j is the characteristic function of $\text{supp } \mu_j$, $j = 1, 2, 3$. The matrix X is piecewise constant; for ease of notation we suppress its z -dependence in the following computations.

Since $\vec{\Phi}'$ is an eigenvector of A corresponding to the eigenvalue $3/2$ (see Section 1.2), we learn from (2.20) and (2.21)

$$X M_1 B\vec{C}_-(z) - X M_2 B\vec{C}_+(z) = -X\vec{\Phi}'(z), \quad z \in \Sigma. \quad (2.22)$$

An immediate consequence of this identity, combined with Sokhotsky–Plemelj's formula $\vec{C}_+ = 2\pi i \vec{\mu}' + \vec{C}_-$, is

$$X M_1 B\vec{C}_+(z) - X M_2 B\vec{C}_-(z) = -X\vec{\Phi}'(z) + 2\pi i X(M_1 + M_2)B\vec{\mu}'(z), \quad z \in \Sigma, \quad (2.23)$$

where as usual, $\vec{\mu} = (\mu_1, \mu_2, \mu_3)^T$, etcetera.

Multiplying (2.19) by X and using (2.22), we conclude that

$$2X A \vec{H}(z) = X(M_1 B \vec{g}_- - M_2 B \vec{g}_+)(z) + h(z) X \vec{\Phi}'(z), \quad z \in \Sigma.$$

Since

$$(M_1 B)^T = -M_2 B, \quad X^T = X,$$

we can write it equivalently as

$$2 \left(\left(A \vec{H} \right)^T X \right)(z) = (\vec{g}_+^T M_1 B X - \vec{g}_-^T M_2 B X)(z) + h(z) \left(\vec{\Phi}'(z) \right)^T X, \quad z \in \Sigma. \quad (2.24)$$

Recall that $\Sigma \subset \Xi_\alpha$ was arbitrary, so we can integrate this formula along $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$.

By (2.17), for $j \in \{1, 2, 3\}$ we get

$$\int g_{j\pm}(y) d\mu_k(y) = \pm \pi i \int h(y) \mu'_j(y) d\mu_k(y) + \int g_j(y) d\mu_k(y), \quad (2.25)$$

where we recall that if g_j is discontinuous across $\text{supp } \mu_k$, we understand the second integral in the right hand side above in terms of its principal value. By assumption,

$\text{supp } h$ does not contain double poles of the coefficient D , and we use the condition (2.16) to get the behavior

$$h(y)\mu'_j(y) = \mathcal{O}(|y-p|^{-\nu}), \quad \text{as } y \rightarrow p \quad (2.26)$$

for some $\nu = \nu(p) \leq 1/3$ and any endpoint p of the support of μ_k . In particular, this implies that $h\mu'_j$ and g_j are μ_k -integrable, so the integrals in the right-hand side of (2.25) are convergent. Moreover, the behavior (2.26) allows us to interchange the order of integration for the second integral in the right-hand side of (2.25) (see [62, Equation (20), page 25]) in order to get

$$\int g_j(y)d\mu_k(y) = - \int h(y)C^{\mu_k}(y)d\mu_j(y).$$

Again by Sokhotsky–Plemelj’s formula,

$$\int h(y)C^{\mu_k}(y)d\mu_j(y) = \mp \pi i \int h(y)\mu'_k(y)d\mu_j(y) + \int h(y)C^{\mu_k}_{\pm}(y)d\mu_j(y),$$

and using it in (2.25), we finally conclude that

$$\int g_{j\pm}(y)d\mu_k(y) = \pm 2\pi i \int \mu'_j(y)\mu'_k(y)h(y)dy - \int h(y)\mu'_j(y)C^{\mu_k}_{\pm}(y)dy.$$

A direct consequence of this formula is that for any 3×3 constant matrix \mathcal{M} ,

$$\int \vec{g}_{\pm}^T \mathcal{M} d\vec{\mu} = \pm 2\pi i \int h(y)(\vec{\mu}'(y))^T \mathcal{M} \vec{\mu}'(y)dy - \int h(y)(\vec{\mu}'(y))^T \mathcal{M} \vec{C}_{\pm}(y)dy.$$

In particular, we have

$$\begin{aligned} \int (\vec{g}_+^T M_1 B - \vec{g}_-^T M_2 B) d\vec{\mu} &= 2\pi i \int h(y)(\vec{\mu}'(y))^T (M_1 + M_2) B \vec{\mu}'(y)dy \\ &\quad - \int h(y)(\vec{\mu}'(y))^T (M_1 B \vec{C}_+(y) - M_2 B \vec{C}_-(y)) dy \\ &= 2\pi i \int h(y)(\vec{\mu}'(y))^T X(y) (M_1 + M_2) B \vec{\mu}'(y)dy \\ &\quad - \int h(y)(\vec{\mu}'(y))^T X(y) (M_1 B \vec{C}_+(y) - M_2 B \vec{C}_-(y)) dy \\ &= \int h(y)(\vec{\mu}'(y))^T \vec{\Phi}'(y)dy, \end{aligned}$$

where we have used the trivial identity $\vec{\mu}'^T X = \vec{\mu}'^T$ for the last two equalities, and also (2.23) for the last equality. Thus, integrating (2.24) with respect to $d\vec{\mu}$, observing that $X d\vec{\mu} = d\vec{\mu}$ and applying this last identity, we arrive at

$$\int (\vec{A} \vec{H})^T(z) d\vec{\mu}(z) = \int h(z) (\vec{\Phi}'(z))^T d\vec{\mu}(z),$$

which is a compact form of writing the condition $\mathcal{D}_h(\vec{\mu}) = 0$. \square

Proof of Theorem 1.8. For a vector-valued function \vec{f} we understand by

$$\int \vec{f}(z)dz$$

the term by term integration, and we denote by $\vec{\mu}' = (\mu'_1, \mu'_2, \mu'_3)$ the vector whose components are the densities of μ_1, μ_2, μ_3 (with respect to the complex line element)

along the arcs of Ξ_α . Because the supports of the components of $\vec{\mu}$ are made of analytic arcs, and $\vec{\mu}'$ has analytic components, we conclude that $\vec{U} = (U^{\mu_1}, U^{\mu_2}, U^{\mu_3})^T$ is continuous across the arcs of Ξ_α .

Integrating (1.11) and using (1.9) we get

$$\operatorname{Re} \int^z \vec{\xi}(y) dy = \operatorname{Re} \left(\frac{1}{3} B \vec{\Phi}(z) \right) + B \vec{U}(z) - \vec{d},$$

where \vec{d} is a real constant vector which only depends on the connected component of $\mathbb{C} \setminus \operatorname{supp} \mu_1 \cup \operatorname{supp} \mu_2 \cup \operatorname{supp} \mu_3$.

Let Σ be an open analytic arc from Ξ_α , not containing any branch point of (1.13). Reasoning as for (2.18), we have

$$2A\vec{U}(z) = (M_1 - M_2) B \vec{U}(z) = M_1 B \vec{U}_-(z) - M_2 B \vec{U}_+(z), \quad z \in \Sigma.$$

Thus, for $z \in \Sigma$ and \vec{d}_\pm the vector of constants for the component on the \pm -side of Σ ,

$$\begin{aligned} 2A\vec{U}(z) &= M_1 \left(\operatorname{Re} \int^z \vec{\xi}_-(y) dy - \operatorname{Re} \left(\frac{1}{3} B \vec{\Phi}(z) \right) + \vec{d}_- \right) \\ &\quad - M_2 \left(\operatorname{Re} \int^z \vec{\xi}_+(y) dy - \operatorname{Re} \left(\frac{1}{3} B \vec{\Phi}(z) \right) + \vec{d}_+ \right) \\ &= \operatorname{Re} \int^z \left(M_1 \vec{\xi}_-(y) - M_2 \vec{\xi}_+(y) \right) dy - \frac{2}{3} \operatorname{Re} A \vec{\Phi}(z) + M_1 \vec{d}_- - M_2 \vec{d}_+. \end{aligned}$$

Recalling that $A\vec{\Phi}' = (3/2)\vec{\Phi}'$ this last equation implies

$$2A\vec{U}(z) = \operatorname{Re} \int^z \left(M_1 \vec{\xi}_-(y) - M_2 \vec{\xi}_+(y) \right) dy - \operatorname{Re} \vec{\Phi}(z) + M_1 \vec{d}_- - M_2 \vec{d}_+, \quad z \in \Sigma.$$

If $\Sigma \subset \operatorname{supp} \mu_{4-j}$, then ξ_j continuous across Σ , and (1.25) tells us that the $(4-j)$ -th entry of $M_1 \xi_- - M_2 \xi_+$ vanishes, which in turn yields (1.26) (with j replaced by $4-j$).

Note that due to (1.26) the S -property on $\operatorname{supp} \mu_{4-j}$ is equivalent to

$$\left[\frac{\partial}{\partial z} \left(\sum_{k=1}^3 a_{4-j,k} U^{\mu_k} + \frac{\phi_{4-j}}{2} \right) \right]_+ = - \left[\frac{\partial}{\partial z} \left(\sum_{k=1}^3 a_{4-j,k} U^{\mu_k} + \frac{\phi_{4-j}}{2} \right) \right]_-, \quad (2.27)$$

Equation (1.11) means that for $z \notin \Xi_\alpha$,

$$2 \frac{\partial}{\partial z} \left(A \vec{U}(z) + \frac{1}{2} \vec{\phi}(z) \right) = A \vec{C}(z) + \frac{1}{2} \vec{\Phi}'(z) = \frac{1}{2} B^T \vec{\xi}(z),$$

where we also used the identity involving A and B in (1.10) and $A\vec{\Phi}' = (3/2)\vec{\Phi}'$. As before, the boundary conditions (1.25) imply that the entry $(4-j)$ of $B^T(\vec{\xi}_- + \vec{\xi}_+)$ vanishes, and the equation above then implies (2.27). \square

Proof of Theorem 1.9. We denote the restriction of Q to the sheet \mathcal{R}_j by Q_j . In cyclic notation mod 3, Q_j^2 can be expressed as

$$Q_j^2(z) = (\xi_{j-1}(z) - \xi_{j+1}(z))^2, \quad z \in \mathcal{R}_j, \quad j = 1, 2, 3.$$

Clearly Q_j^2 is meromorphic in $\overline{\mathbb{C}} \setminus (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3)$. We first prove that each Q_j^2 is meromorphic on the sheet \mathcal{R}_j . If $\Gamma \subset \text{supp } \mu_3$, then due to (1.6) the function ξ_1 is analytic across Γ , hence

$$\xi_{2\pm}(s) = \xi_{3\mp}(s), \quad s \in \Gamma,$$

implying that

$$Q_{1+}^2(s) = (\xi_{2+}(s) - \xi_{3+}(s))^2 = (\xi_{3-}(s) - \xi_{2-}(s))^2 = Q_{1-}^2(s), \quad s \in \Gamma,$$

so Q_1^2 is analytic across Γ and, as a consequence, it follows that Q_1^2 is meromorphic on the whole sheet \mathcal{R}_1 . Similarly we prove Q_j^2 is meromorphic on \mathcal{R}_j , $j = 2, 3$.

We now show that Q^2 is globally defined, that is, the function Q_{j+1}^2 is the analytic continuation of Q_j^2 from \mathcal{R}_j to \mathcal{R}_{j+1} , $j = 1, 2, 3$.

Consider an arc $\Gamma \subset \Xi_\alpha$. For the construction of \mathcal{R} we know that Γ connects exactly two sheets, which we assume to be \mathcal{R}_1 and \mathcal{R}_2 , the remaining cases are analogous. Then $\Gamma \subset \text{supp } \mu_1 \setminus (\text{supp } \mu_2 \cup \text{supp } \mu_3)$, so the function ξ_3 is analytic across Γ , hence $\xi_{1\pm} = \xi_{2\mp}$ on Γ and

$$Q_{1\pm}^2(s) = (\xi_{2\pm}(s) - \xi_3(s))^2 = (\xi_{1\mp}(s) - \xi_3(s))^2 = Q_{2\mp}^2(s),$$

so Q_2^2 is the analytic continuation of Q_1^2 to \mathcal{R}_2 across Γ .

Thus, Q^2 is meromorphic along any arc connecting the sheets. It is clear that the functions ξ_1, ξ_2, ξ_3 , and hence Q^2 , can only be unbounded at the points at ∞ and also at the poles of the coefficient D in (1.13). If p is such a pole, then the three functions ξ_1, ξ_2, ξ_3 are branched at p , and the local coordinate of \mathcal{R} at $z = p$ is of the form $z = p + u^3$, $u \in \mathbb{C}$, and hence

$$dz^2 = 9u^4 du^2. \quad (2.28)$$

Since p is at most a double pole of D , it follows from the identities

$$D(z) = -\xi_j(z)^3 + R(z)\xi_j(z), \quad j = 1, 2, 3,$$

that the functions ξ_1, ξ_2, ξ_3 all blow up with the same order $\nu/3$, $\nu \in \{1, 2\}$, that is

$$\xi_j(z) = \frac{\kappa_j}{z^{\nu/3}}(1 + o(1)), \quad \text{as } z \rightarrow p, \quad j = 1, 2, 3,$$

for some nonzero constants $\kappa_1, \kappa_2, \kappa_3$, so

$$Q^2(z) = \frac{c}{z^{2\nu/3}}(1 + o(1)) = \frac{c}{u^{2\nu}}(1 + o(1)), \quad \text{as } z \rightarrow p,$$

for some constant c . Combining with (2.28), we get

$$-Q(z)^2 dz^2 = \tilde{c} u^{4-2\nu} (1 + o(1)) du^2,$$

for some constant \tilde{c} . Since $\nu \leq 2$, we get that ϖ has to remain bounded near p , and hence ϖ can only have poles at the points at ∞ .

Let $\Sigma \subset \Xi_\alpha \cap \text{supp } \mu_1$ be an open arc. Due to assumption (1.6), function ξ_3 is analytic across Σ . Now the definition of function ξ_1 in (1.12), the Sokhotsky-Plemelj's formula and (1.25) yield that for $s \in \Sigma$,

$$d\mu_1(s) = d\mu_1(s) + d\mu_2(s) = \frac{1}{2\pi i} (\xi_{1+}(s) - \xi_{1-}(s)) ds = \frac{1}{2\pi i} (\xi_{1+}(s) - \xi_{2+}(s)) ds.$$

Analogously, if $\Sigma \subset \Xi_\alpha \cap \text{supp } \mu_2$, then ξ_2 is analytic across Σ , and for $s \in \Sigma$,

$$d\mu_2(s) = d\mu_1(s) + d\mu_2(s) = \frac{1}{2\pi i} (\xi_{1+}(s) - \xi_{1-}(s)) ds = \frac{1}{2\pi i} (\xi_{1+}(s) - \xi_{3+}(s)) ds.$$

Finally, if $\Sigma \subset \Xi_\alpha \cap \text{supp } \mu_3$, then ξ_1 is analytic across Σ , and using again the definition of ξ_3 in (1.12) and the Sokhotsky-Plemelj's formula, we get that for $s \in \Sigma$,

$$d\mu_3(s) = d\mu_3(s) - d\mu_1(s) = \frac{1}{2\pi i} (\xi_{3+}(s) - \xi_{3-}(s)) ds = \frac{1}{2\pi i} (\xi_{3+}(s) - \xi_{2+}(s)) ds.$$

This shows that the measures μ_j must be supported on arcs of trajectories of ϖ . \square

3. THE CUBIC CASE

In the following two sections we deal with the cubic case, and describe a one-parameter family of critical vector-valued measures for the energy functional (1.1)–(1.2) and for the choice (1.30), so that the external fields are given by (1.31). Although the α -critical measures were defined for $\alpha \in [0, 1]$, here we restrict our attention to $\alpha \in [0, 1/2]$. As it will follow from our analysis below, as $\alpha \nearrow 1/2$, the support of the component μ_3 of the α -critical measure becomes unbounded. Furthermore, our original motivation was the asymptotic analysis of the multiple orthogonal polynomials (1.36), for which the case of $\alpha \in [1/2, 1]$ can be easily reduced to $\alpha \in [0, 1/2]$ by an appropriate rotation of the plane. Thus, the selection of this interval for α is natural in the present situation, although our method carries over without any special difficulty to the whole range of α .

3.1. The spectral curve. Recall that by Theorem 1.2 the shifted resolvents, defined in (1.12), of any critical vector-valued measure $\vec{\mu} = (\mu_1, \mu_2, \mu_3)$ satisfy the algebraic equation (“spectral curve”) (1.13). As a first step, we deduce the expressions (1.32)–(1.33) for its coefficients.

For the potentials as in (1.31), the coefficient R in (1.14) reduces to

$$R(z) = 3z^4 - 3z - c, \quad (3.1)$$

where at this moment the constant c is given in the form

$$c = \int x(d\mu_1(x) + d\mu_2(x)).$$

Since $D(z) = -\xi_1^3 + R(z)\xi_1 = \xi_1\xi_2(\xi_1 + \xi_2)$, comparing the expansions of both expressions at ∞ and using (1.4), we further get that

$$C^{\mu_1}(z) + C^{\mu_2}(z) = -\frac{1}{z} - \frac{c}{3z^2} + \frac{1 - \alpha(1 - \alpha)}{3z^4} + \mathcal{O}\left(\frac{1}{z^5}\right), \quad z \rightarrow \infty,$$

and

$$D(z) = -2z^6 + 3z^3 + cz^2 - 3\tau, \quad \tau := \alpha(1 - \alpha). \quad (3.2)$$

Observe that $\tau = \alpha(1 - \alpha)$ is an equivalent parametrization that gives a bijection between the interval $\alpha \in [0, \frac{1}{2}]$ and $\tau \in [0, \frac{1}{4}]$. As it was mentioned in Section 1.4, the extremal cases, $\tau = 0$ and $\tau = \frac{1}{4}$, were studied in [27] and [35], respectively, in their connection to the multiple orthogonal polynomials (1.36), for the choices $n = 0$ and $n = m$, respectively, and the family of critical measures depending on $\alpha \in (0, 1/2)$ are also relevant to the asymptotic analysis of these polynomials for general n, m . From a different perspective, we are studying a continuous deformation of the critical measures, interpolating the extremal cases of $\tau = 0$ and $\tau = \frac{1}{4}$.

The discriminant of (1.13) with respect to the variable ξ is

$$\begin{aligned} \text{Discr}(z) &= 4R^3(z) - 27D^2(z) \\ &= (81 - 324\tau)z^6 + 54cz^5 + 9c^2z^4 + 54(9\tau - 2)z^3 \\ &\quad + 54c(3\tau - 2)z^2 - 36c^2z - 4c^3 - 243\tau^2. \end{aligned}$$

Since $\text{Discr}(z)$ is a polynomial in z of degree 6, if the Riemann surface of (1.13) has genus 0 then, by the Riemann-Hurwitz theorem, $\text{Discr}(z)$ must have a multiple root. In particular, the discriminant Discr_1 of $\text{Discr}(z)$ must vanish. A cumbersome but straightforward calculation (that can be carried out with the aid of a symbolic algebra software such as Mathematica) shows that

$$\text{Discr}_1 = a\tau p_1(c)p_2(c)^3$$

for a nonzero real constant a and

$$\begin{aligned} p_1(c) &= 64c^3 + 243(1 - 4\tau)^2, \\ p_2(c) &= c^6 - 486c^3\tau(1 + \tau) + 2187\tau(3\tau - 1)^3, \end{aligned} \tag{3.3}$$

so we expect c to be a root of either p_1 or p_2 . For the case $\tau = \frac{1}{4}$ studied in [35], the algebraic equation (1.13) reduces to

$$\xi^3 - (3z^4 - 3z)\xi - 2z^6 + 3z^3 - \frac{3}{4} = 0,$$

showing that $c = 0$, which is a root of p_1 . By continuity, we expect c to be a root of p_1 for every choice of $\tau \in (0, \frac{1}{4})$, concluding that

$$c^3 = -\frac{243}{64}(1 - 4\tau)^2,$$

or

$$c = -\left(\frac{243}{64}(1 - 4\tau)^2\right)^{\frac{1}{3}}. \tag{3.4}$$

which is the same as (1.33).

Obviously, c defined by (3.4) can take three possible values; for the rest of the paper we choose c to be real (and thus, negative), so that the algebraic equation (1.13) is also real. It should be pointed out that a different choice of c would lead to an algebraic equation corresponding to another triplet of critical measures $\vec{\mu}$, that can be obtained from the original one by rotation by $\pm 2\pi/3$.

For the choice of c in (3.4), we can rewrite $\text{Discr}(z)$ as

$$\text{Discr}(z) = \frac{243}{256}q_1(z)q_2(z)^2, \tag{3.5}$$

where q_1, q_2 are given by

$$\begin{aligned} q_1(z) &= \frac{256}{3^{5/3}}(1 - 4\tau)^{1/3}z^4 + \frac{128}{9}z^3 + \frac{16}{3^{1/3}}(1 - 4\tau)^{2/3}z^2 \\ &\quad - 3^{1/3}32(1 - 4\tau)^{1/3}z + 16(1 - 8\tau), \end{aligned} \tag{3.6}$$

$$q_2(z) = 3^{1/3}(1 - 4\tau)^{1/3}z - 1. \tag{3.7}$$

For $\tau < \frac{1}{4}$ the discriminant of q_1 is

$$\text{const} \times \tau(1 - \tau)^2(27 - 100\tau)^3,$$

which never vanishes, and we conclude that q_1 has always four distinct roots. For the choice $\tau = \frac{1}{8}$, q_1 has two complex conjugate roots and two real roots, so this also holds for any τ .

The resultant of q_1 and q_2 is

$$\text{const} \times (1 - 12\tau)^4 (1 - 4\tau)^{2/3},$$

so that in the interval $0 < \tau < 1/4$, the polynomials q_1, q_2 share a root only when $\tau = 1/12$. It then follows that for $\tau \neq 1/12$ the roots of q_1 are branch points of multiplicity 2 of (1.13).

If the double zero b_* of Discr (that is, the zero of q_2) is a branch point of (1.13), then its multiplicity (as a branch point) has to be three. This means that the three solutions to (1.13) should coincide for $z = b_*$, so (1.13) has to share a root with its second ξ -derivative, and hence $\xi_j(b_*) = 0$, $j = 1, 2, 3$. Plugging this back into (1.13) we see that $D(b_*) = 0$. But

$$D(b_*) = -\frac{(12\tau - 1)^3}{36(1 - 4\tau)^2},$$

so $D(b_*) = 0$ only for $\tau = 1/12$. Hence for $\tau \neq 1/12$ the point b_* is a regular point of (3.9), that is, the algebraic equation (1.13) is not branched at $z = b_*$. We already observed that the simple zeros of Discr are always branch points of multiplicity 2, so the Riemann-Hurwitz formula says that for $\tau \neq 1/12$ the associated algebraic equation has genus 0. Continuity with respect to τ assures that the genus cannot increase for $\tau = 1/12$, and hence the genus is also zero for this value.

The discussion above can be summarized in the following proposition:

Proposition 3.1. *The algebraic equation (1.13) with coefficients given by (3.1), (3.2) and (3.4), has four branch points (the zeros of the polynomial q_1 in (3.6)) and a double point (the zero of q_2 in (3.7)). Two of the branch points are real and the other two form a complex conjugate pair. For $\tau \neq 1/12$ all these points are distinct, while when $\tau = 1/12$, the double point and one of the real branch points of (1.13) coalesce. The associated Riemann surface has always genus 0.*

Although the choice of c in (3.4), as a root of p_1 instead of p_2 , was mostly motivated by the construction of a continuous one-parameter family of critical measures $\bar{\mu}$ interpolating the extremal cases $\alpha = 0$ and $\alpha = 1/2$ studied in [27, 35], formulas (3.1), (3.2) and (3.4) can also be explained in terms of the genus 0 ansatz stated in Theorem 1.12 that we prove next.

Proof of the second part of Theorem 1.12. We already noted that if the Riemann surface has genus 0, then c must be a root of at least one of the polynomials p_1 and p_2 in (3.3), and Proposition 3.1 assures that the roots of p_1 give rise to genus 0 Riemann surfaces. It thus suffices to show that zeros of p_2 give rise to a surface of genus 1.

Consider c to be a root of the polynomial p_2 . As before, c can be assumed to be real, the remaining nonreal choices of c as a root of p_2 can be reduced to the real ones with suitable change of variables. Hence,

$$c = 3(9\tau + 9\tau^2 + \sqrt{3}(1 + 9\tau)\tau^{1/2})^{1/3} \quad \text{or} \quad c = 3(9\tau + 9\tau^2 - \sqrt{3}(1 + 9\tau)\tau^{1/2})^{1/3}. \quad (3.8)$$

Making the change of variables $\tau = u^6/3$, this can be written as

$$c = c(u) = 3u(u + 1)(u^2 - u + 1), \quad 0 < |u| \leq \left(\frac{3}{4}\right)^{1/6},$$

where $u > 0$ corresponds to the first value of c in (3.8) and $u < 0$ corresponds to the second choice of c in (3.8).

With this identification, the discriminant $\text{Discr}(z)$ simplifies to

$$\text{Discr}(z) = -27(u+z)^2 q(z),$$

where

$$\begin{aligned} q(z) = & z^4(4u^6 - 3) - z^3(8u^7 + 6u^4) + z^2(9u^8 + 6u^5) \\ & - z(10u^9 + 12u^6 - 4) + 5u^{10} + 12u^7 + 12u^4 + 4u. \end{aligned}$$

The discriminant of q is given by

$$-6912(u^3 + 1)^6 (10u^3 + 9) (3u^6 + 1)^3 \neq 0 \quad \text{for } |u| \in (0, (3/4)^{1/6}).$$

Moreover,

$$q(-u) = 9u^4 (2u^3 + 1)^2,$$

hence $q(-u) \neq 0$, unless $u = -1/2^{1/3}$, which corresponds to $\tau = 1/12$. But for this latter choice, the corresponding value of c in (3.8) coincides with the value of c in (3.4), hence the coefficients R and D are given by (1.32).

Thus, let $u \neq -1/2^{1/3}$. Then $\text{Discr}(z)$ has four simple roots, which have to be branch points of the equation, and one double root $z = -u$. At this double root, the algebraic equation (1.13) simply reduces to $\xi^3 = 0$, so its three solutions coincide. This is only compatible with the fact that $\text{Discr}(z)$ has a double root at $z = -u$ if the three solutions are branched at this point. Hence we have four branch points of multiplicity 2 and one branch point of multiplicity 3, and the Riemann-Hurwitz formula gives us that the Riemann surface has genus 1. \square

The first part of Theorem 1.12, claiming existence of α -critical measures, will be given by Corollary 3.6.

Remark 3.2. Obviously the quadratic differential (1.28) still makes sense if c is a root of p_2 as in the proof above. Numerical experiments performed to compute its critical graph indicate that the corresponding algebraic equation should also give rise to α -critical measures as in Theorem 1.3.

3.2. Equilibrium problem from the spectral curve. In this section, starting from the algebraic equation

$$\xi^3 - R(z)\xi + D(z) = 0, \tag{3.9}$$

where it is assumed that the coefficients R and D are given by (3.1)–(3.4), we find a vector of measures $\vec{\mu} = (\mu_1, \mu_2, \mu_3) \in \mathcal{M}_\alpha$ for which the respective ξ -functions in (1.12) satisfy (3.9), and consequently we prove Theorems 1.14 and 1.16. A central object for this analysis is the associated quadratic differential (1.28) (see also (3.11) *et seq.* below) and its critical graph, whose description is postponed to Section 4.

According to Proposition 3.1, the spectral curve (3.9) has two real branch points $a_1 < b_1$, two non real branch points $b_2 = \bar{a}_2$, with $\text{Re } b_2 > 0$, and a double point

$$b_* = \frac{1}{(3(1-4\tau))^{1/3}} > 0 \quad \text{for } \tau \in (0, 1/4).$$

Since

$$q_1(b_*) = \frac{32(1-12\tau)^2}{9(1-4\tau)} \geq 0 \quad \text{for } \tau \in (0, 1/4),$$

we easily conclude that $b_* \geq b_1$ for $\tau \in (0, 1/4)$, with $b_* = b_1$ if and only if $\tau = 1/12$.

Moreover, equation (3.9) defines a three-sheeted Riemann surface \mathcal{R} of genus 0 for every value of the parameter $\tau \in (0, 1/4)$. In contrast to (1.27), where the cuts for the Riemann surface are defined in terms of the supports of the critical measures, here the cuts that split \mathcal{R} into the sheets $\mathcal{R}_1, \mathcal{R}_2, \mathcal{R}_3$ can be chosen in a somewhat arbitrary way, as long as they connect the branch points. In our case, the sheet structure will be given by two cuts: the interval of the real line connecting the real branch points a_1, b_1 of \mathcal{R} and a Jordan arc (to be defined precisely below) that joins the complex-conjugate branch points a_2, b_2 . This arc intersects \mathbb{R} in a unique point a_* , whose location depends on the value of the parameter τ . Namely, there is a certain critical value $\tau_c \approx 0.19$, defined later, such that $a_* \in (a_1, b_1)$ if and only if $\tau_c < \tau < 1/4$ (what we call the *supercritical regime*), see Figure 2. It is important to stress here that $\tau_c > 1/12$.

Obviously, it ultimately follows from our arguments that this cut structure coincides with (1.27) for the critical measure given by Theorem 1.14.

Using these cuts we define three oriented sets, Δ_1, Δ_2 and Δ_3 on \mathbb{C} . Namely, Δ_2 is always the projection onto \mathbb{C} of the cut joining the complex-conjugate branch points a_2, b_2 and oriented from a_2 to b_2 . In the subcritical regime ($0 < \tau < \tau_c$) we denote by $\Delta_1 = [a_1, b_1] \subset \mathbb{R}$ oriented from a_1 to b_1 , and set $\Delta_3 = \emptyset$. In the supercritical regime ($\tau_c < \tau < 1/4$), $\Delta_1 = [a_*, b_1] \subset \mathbb{R}$ and $\Delta_3 = [a_1, a_*] \subset \mathbb{R}$, both with the natural orientation. The orientation induces the left (denoted by the subscript “+”) and right (with the subscript “−”) boundary values of functions defined on \mathbb{C} or \mathcal{R} .

We define

$$\mathcal{R}_1 = \overline{\mathbb{C}} \setminus (\Delta_1 \cup \Delta_2), \quad \mathcal{R}_2 = \overline{\mathbb{C}} \setminus (\Delta_1 \cup \Delta_3), \quad \mathcal{R}_3 = \overline{\mathbb{C}} \setminus (\Delta_2 \cup \Delta_3),$$

and build the Riemann surface $\mathcal{R} = \mathcal{R}_1 \cup \mathcal{R}_2 \cup \mathcal{R}_3$, associated to the algebraic equation (3.9), connecting each pair of sheets \mathcal{R}_j crosswise across the cuts as follows: Δ_1 connects \mathcal{R}_1 with \mathcal{R}_2 , Δ_2 connects \mathcal{R}_1 with \mathcal{R}_3 , and Δ_3 connects \mathcal{R}_2 with \mathcal{R}_3 (this last condition is clearly non-trivial only in the supercritical regime, when $\tau_c < \tau < 1/4$), see again Figure 2.

As we will see soon, the cuts $\Delta_1, \Delta_2, \Delta_3$ can be specified in such a way that we will be able to define positive measures living on Δ_j ’s, and the construction of the sheet structure will coincide with (1.27).

The three solutions ξ_1, ξ_2, ξ_3 of (3.9) are enumerated according to their asymptotic expansion at infinity:

$$\begin{aligned} \xi_1(z) &= 2z^2 - \frac{1}{z} + \mathcal{O}(z^{-2}), \\ \xi_2(z) &= -z^2 + \frac{\alpha}{z} + \mathcal{O}(z^{-2}), \\ \xi_3(z) &= -z^2 + \frac{1-\alpha}{z} + \mathcal{O}(z^{-2}). \end{aligned} \tag{3.10}$$

As before, the functions ξ_1, ξ_2, ξ_3 are regarded as branches of the same meromorphic function $\xi : \mathcal{R} \rightarrow \mathbb{C}$, defined by (3.9). As it is rigorously given by Proposition 4.3 below, the function ξ_j is the restriction of ξ to the sheet \mathcal{R}_j of the Riemann surface \mathcal{R} .

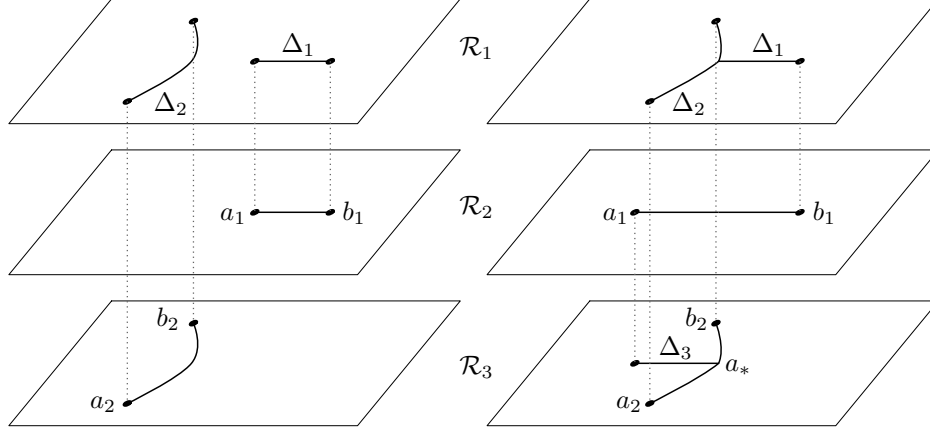


FIGURE 2. Sheet structure in the *precritical* $0 < \tau < \tau_c$ (left) and *supercritical* $\tau_c < \tau < 1/4$ (right) regimes.

Due to the explicit geometric description of the sets Δ_j 's and sheets \mathcal{R}_j 's, it is straightforward to check that with

$$Q(z) = \begin{cases} \xi_2(z) - \xi_3(z), & \text{on } \mathcal{R}_1, \\ \xi_1(z) - \xi_3(z), & \text{on } \mathcal{R}_2, \\ \xi_1(z) - \xi_2(z), & \text{on } \mathcal{R}_3, \end{cases} \quad (3.11)$$

Q^2 extends as a single-valued meromorphic function on \mathcal{R} , and

$$\varpi = -Q^2(z) dz^2 \quad (3.12)$$

is the corresponding rational quadratic differential on \mathcal{R} . The details are also carried out in the general setting of Theorem 1.9 in its proof.

One of the main outcomes of the discussion in Section 4 (see Theorem 4.4 and Remark 4.11) is that we can choose the cut Δ_2 connecting a_2, b_2 to coincide with the trajectory along which

$$\operatorname{Re} \int^z (\xi_{1+}(s) - \xi_{3+}(s)) ds = \text{const}, \quad z \in \Delta_2. \quad (3.13)$$

Furthermore, Δ_2 can be extended to an analytic arc Γ from $e^{-\frac{2\pi i}{3}}\infty$ to $e^{\frac{2\pi i}{3}}\infty$ in such a way that

$$\operatorname{Im} \int^z (\xi_{1+}(s) - \xi_{3+}(s)) ds = \text{const}, \quad z \in \Gamma \setminus \Delta_2, \quad (3.14)$$

and

$$\xi_1(z) - \xi_3(z) \neq 0, \quad z \in \Gamma \setminus \{a_2, b_2\}, \quad (3.15)$$

see Figure 3.

Further,

$$\operatorname{Re} \int^z (\xi_{1+}(s) - \xi_{2+}(s)) ds = \text{const}, \quad z \in \Delta_1, \quad (3.16)$$

and

$$\xi_{1+}(z) - \xi_{2+}(z) \neq 0, \quad z \in \Delta_1 \setminus \{a_1, b_1\}, \quad (3.17)$$

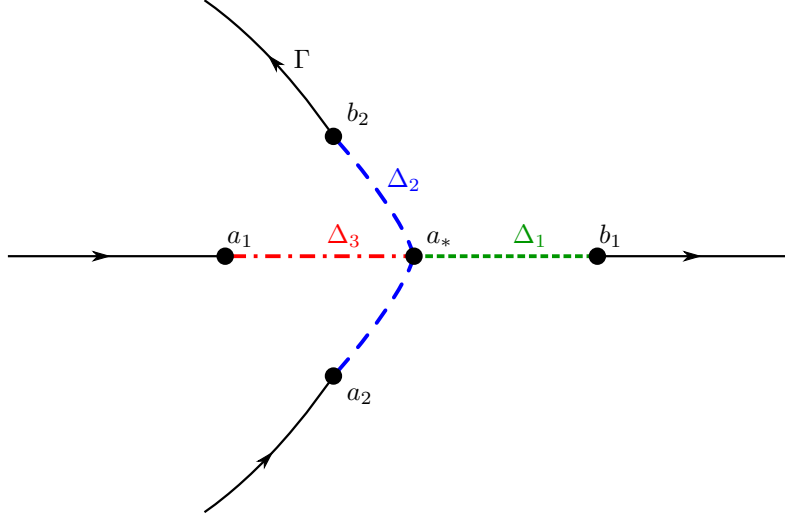


FIGURE 3. Pictorial representation of the sets Δ_1 , Δ_2 , Δ_3 and Γ in the supercritical regime.

and in the supercritical regime $\tau_c < \tau < 1/4$, also

$$\operatorname{Re} \int^z (\xi_{2+}(s) - \xi_{3+}(s)) ds = \text{const}, \quad z \in \Delta_3, \quad (3.18)$$

and

$$\xi_{2+}(s) - \xi_{3+}(s) \neq 0, \quad z \in \Delta_3 \setminus \{a_1, a_*\},$$

for which we refer to Proposition 4.3.

Finally, it also holds

$$\begin{aligned} \xi_1(z) - \xi_2(z) &> 0, \quad z > b_1, \\ \xi_2(z) - \xi_3(z) &> 0, \quad z < \min\{a_1, a_*\}, \end{aligned} \quad (3.19)$$

and in the precritical regime $\tau < \tau_c$,

$$\xi_1(z) - \xi_2(z) < 0, \quad z \in (a_*, a_1), \quad (3.20)$$

for which again we refer to Proposition 4.3.

We define measures μ_1 , μ_2 , μ_3 on Δ_1 , Δ_2 , Δ_3 , respectively, through the formulas

$$\begin{aligned} d\mu_1(s) &= \frac{1}{2\pi i} (\xi_{1+}(s) - \xi_{2+}(s)) ds, \quad s \in \Delta_1, \\ d\mu_2(s) &= \frac{1}{2\pi i} (\xi_{1+}(s) - \xi_{3+}(s)) ds, \quad s \in \Delta_2, \\ \text{and} \end{aligned} \quad (3.21)$$

$$d\mu_3(s) = \begin{cases} 0, & \text{if } 0 < \tau < \tau_c, \\ \frac{1}{2\pi i} (\xi_{3+}(s) - \xi_{2+}(s)) ds, & s \in \Delta_3, \text{ if } \tau_c < \tau < 1/4, \end{cases}$$

where ds denotes the complex line element on the respective arc. Due to their construction, the supports of the measures μ_1 , μ_2 , μ_3 satisfy the claims of Theorem 1.14 (see also Theorem 4.4).

Remark 3.3. Using the definition in (3.11) we can describe formulas (3.21) saying that we build μ_j from the values of Q on the sheet \mathcal{R}_{4-j} , $j = 1, 2, 3$.

Proposition 3.4. *Expressions (3.21) define positive measures μ_1, μ_2, μ_3 .*

Proof. We prove the statement for the supercritical regime $\tau_c < \tau < 1/4$; the precritical regime, somewhat simpler, can be analyzed similarly.

From (3.13) and (3.16) it follows that the measures μ_1, μ_2, μ_3 are real. Moreover, the densities of μ_1, μ_3 with respect to the complex line element ds are continuous and non vanishing in the interior of Δ_1, Δ_3 , so the respective measures do not change sign. As for μ_2 , its density with respect to ds is continuous when restricted to either the upper or the lower half plane, but not across \mathbb{R} .

We start now with μ_1 . For $x \in \Delta_1$, we can write

$$\mu_1([x, b_1]) = -\frac{1}{2\pi i} \psi_{1+}(x), \quad (3.22)$$

where

$$\psi_1(z) = \int_{b_1}^z (\xi_1(s) - \xi_2(s)) ds, \quad z \in \mathbb{C} \setminus ((-\infty, b_1] \cup \Delta_2).$$

By the asymptotic expansion (3.10), we have that

$$\psi_1(z) = z^3 + \mathcal{O}(z), \quad z \rightarrow \infty,$$

and since $\xi_1 - \xi_2$ does not change sign on $[b_1, +\infty)$, see (3.17), we conclude that

$$\psi_1(x) > 0, \quad x > b_1.$$

The function ψ_1 vanishes at the point $z = b_1$ with order $3/2$. Since we already know that ψ_1 should map Δ_1 to $i\mathbb{R}$, from its order of vanishing we further get that the imaginary part of $\psi_{1+}(x)$ is negative when $x \in \Delta_1$ is sufficiently close to the endpoint b_1 . From (3.17), we conclude

$$\psi_{1+}(x) \in i\mathbb{R}_-, \quad x \in \Delta_1.$$

By (3.22), this establishes the positivity of μ_1 .

Furthermore, the density $d\mu_3/ds$ is the analytic continuation of $d\mu_1/ds$ across Δ_2 , hence $d\mu_3$ is also positive.

Finally, in order to establish the sign of μ_2 , consider first $x \in \Delta_2 \cap \mathbb{C}_-$, $\mathbb{C}_- = \{z \in \mathbb{C} : \text{Im}(z) < 0\}$. Denote by $\Delta_2[a_2, x]$ the subarc of Δ_2 from a_2 to x . Then

$$\mu_2(\Delta_2[a_2, x]) = \frac{1}{2\pi i} \psi_{2+}(x), \quad (3.23)$$

where

$$\psi_2(z) = \int_{a_2}^z (\xi_{1+}(s) - \xi_{3+}(s)) ds, \quad z \in \mathbb{C} \setminus ((-\infty, b_1] \cup \Delta_2).$$

From (3.14) and (3.15) we learn that $\text{Re}(\psi_2)$ is monotone on $\Gamma(e^{-\frac{3\pi i}{2}} \infty, a_2)$. The asymptotics

$$\psi_2(z) = z^3 + \mathcal{O}(z), \quad z \rightarrow \infty,$$

shows that

$$\text{Re } \psi_2(z) \rightarrow +\infty, \quad z \rightarrow \infty \text{ along } \Gamma,$$

hence $\text{Re}(\psi_2)$ is strictly decreasing on $\Gamma(-\frac{3\pi i}{2} \infty, a_2)$, $\psi_2(a_2) = 0$. Since ψ_2 vanishes with order $3/2$ on a_2 , we conclude that

$$\psi_{2+}(x) \in i\mathbb{R}_+, \quad x \in \Delta_2 \cap \mathbb{C}_-.$$

Thus, by (3.23) measure μ_2 is positive on $\Delta_2 \cap \mathbb{C}_-$.

The positivity of μ_2 on $\Delta_2 \cap \mathbb{C}_+$ can be obtained by similar arguments or using the real symmetry of the density of μ_2 ; we leave the details to the interested reader.

The proposition is proved. \square

Our next goal is an expression for the Cauchy transforms of combinations of measures μ_1, μ_2, μ_3 :

Proposition 3.5. *The Cauchy transform of the measures μ_1, μ_2, μ_3 defined in (3.21) are related to the ξ -functions in (3.10) through*

$$C^{\mu_1}(z) + C^{\mu_2}(z) + 2z^2 = \xi_1(z), \quad z \in \mathbb{C} \setminus (\Delta_1 \cup \Delta_2), \quad (3.24)$$

$$C^{\mu_1}(z) + C^{\mu_3}(z) + z^2 = -\xi_2(z), \quad z \in \mathbb{C} \setminus (\Delta_1 \cup \Delta_3), \quad (3.25)$$

$$C^{\mu_2}(z) - C^{\mu_3}(z) + z^2 = -\xi_3(z), \quad z \in \mathbb{C} \setminus (\Delta_2 \cup \Delta_3). \quad (3.26)$$

In particular, the total masses of the measures μ_1, μ_2, μ_3 satisfy (1.4), namely

$$|\mu_1| + |\mu_2| = 1, \quad |\mu_1| + |\mu_3| = \alpha, \quad |\mu_2| - |\mu_3| = 1 - \alpha.$$

Proof. This is a straightforward consequence of residues calculations. For instance,

$$\begin{aligned} C^{\mu_1}(z) + C^{\mu_2}(z) &= \frac{1}{2\pi i} \int_{\Delta_1} \frac{\xi_{1+}(s) - \xi_{1-}(s)}{s - z} ds + \frac{1}{2\pi i} \int_{\Delta_2} \frac{\xi_{1+}(s) - \xi_{1-}(s)}{s - z} ds \\ &= \frac{1}{2\pi i} \oint \frac{\xi_1(s)}{s - z} ds \\ &= \text{Res} \left(\frac{\xi_1(s)}{s - z}, s = z \right) + \text{Res} \left(\frac{\xi_1(s)}{s - z}, s = \infty \right) \\ &= -2z^2 + \xi_1(z). \end{aligned}$$

where the closed integral is computed along a contour oriented clockwise which encircles $\Delta_1 \cup \Delta_2$ and does not encircle z .

Analogously,

$$\begin{aligned} C^{\mu_1}(z) + C^{\mu_3}(z) &= -\text{Res} \left(\frac{\xi_2(s)}{s - z}, s = z \right) - \text{Res} \left(\frac{\xi_2(s)}{s - z}, s = \infty \right) \\ &= -z^2 - \xi_2(z), \end{aligned}$$

and taking the difference of these last two equations,

$$C^{\mu_2}(z) - C^{\mu_3}(z) = -z^2 + \xi_1(z) + \xi_2(z) = -z^2 - \xi_3(z).$$

Finally, (1.4) is a direct consequence of the asymptotic expansion (3.10). \square

A combination of Proposition 3.5 with Theorem 1.3 yields the first part of Theorem 1.12, namely

Corollary 3.6. *The vector of measures $\vec{\mu} = (\mu_1, \mu_2, \mu_3) \in \mathcal{M}_\alpha$ defined through (3.21) is α -critical for the potentials (1.31).*

Recalling (1.9), the potential U^μ of a compactly supported signed measure μ for which $\mathbb{C} \setminus \text{supp } \mu$ is connected is the real part of a primitive of C^μ , that is

$$U^\mu(z) = \text{Re} \int^z C^\mu(s) ds + c, \quad z \in \mathbb{C} \setminus \text{supp } \mu, \quad (3.27)$$

where the constant is chosen so as to have

$$\lim_{z \rightarrow \infty} \left(\operatorname{Re} \int^z C^\mu(s) ds + c \right) = 0.$$

Let us apply this to the measures μ_1, μ_2, μ_3 given in (3.21). From (3.24),

$$\int_{b_1}^z (C^{\mu_1}(s) + C^{\mu_2}(s)) ds + \frac{2}{3}z^3 - \frac{2}{3}b_1^3 = \int_{b_1}^z \xi_1(s) ds, \quad z \in \mathbb{C} \setminus (\Delta_1 \cup \Delta_2),$$

and from (3.25),

$$\int_{b_1}^z (C^{\mu_1}(s) + C^{\mu_3}(s)) ds + \frac{1}{3}z^3 - \frac{1}{3}b_1^3 = - \int_{b_1}^z \xi_2(s) ds, \quad z \in \mathbb{C} \setminus (\Delta_1 \cup \Delta_3).$$

Summarizing,

$$\begin{aligned} \int_{b_1}^z (2C^{\mu_1}(s) + C^{\mu_2}(s) + C^{\mu_3}(s)) ds + z^3 - b_1^3 \\ = \int_{b_1}^z (\xi_1(s) - \xi_2(s)) ds, \quad z \in \mathbb{C} \setminus (\Delta_1 \cup \Delta_2 \cup \Delta_3), \end{aligned}$$

where we must use the same paths of integration in the left and in the right hand sides. By (3.27) we then conclude

$$2U^{\mu_1}(z) + U^{\mu_2}(z) + U^{\mu_3}(z) + \phi(z) - l_1 = \operatorname{Re} \int_{b_1}^z (\xi_1(s) - \xi_2(s)) ds, \quad z \in \mathbb{C} \setminus \Delta, \quad (3.28)$$

for some constant l_1 , the external field ϕ given in (1.31) and $\Delta = \Delta_1 \cup \Delta_2 \cup \Delta_3$. Since the ξ_j 's have purely imaginary periods, the right hand side above is well defined regardless the path chosen.

In a completely analogous way, we get

$$\begin{aligned} 2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 &= \operatorname{Re} \int_{a_2}^z (\xi_1(s) - \xi_3(s)) ds, \quad z \in \mathbb{C} \setminus \Delta, \\ 2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - l_3 &= \operatorname{Re} \int_{\min\{a_1, a_*\}}^z (\xi_3(s) - \xi_2(s)) ds, \quad z \in \mathbb{C} \setminus \Delta. \end{aligned} \quad (3.29)$$

Combining (3.28) with (3.16), we conclude that

$$2U^{\mu_1}(z) + U^{\mu_2}(z) + U^{\mu_3}(z) + \phi(z) - l_1 = 0, \quad z \in \Delta_1. \quad (3.30)$$

Moreover, from the first equation in (3.19) we also get

$$2U^{\mu_1}(z) + U^{\mu_2}(z) + U^{\mu_3}(z) + \phi(z) - l_1 > 0, \quad z \in (b_1, +\infty). \quad (3.31)$$

Furthermore, in the precritical case $\tau < \tau_c$ and $z \in [a_*, a_1)$,

$$\begin{aligned} 2U^{\mu_1}(z) + U^{\mu_2}(z) + U^{\mu_3}(z) + \phi(z) - l_1 &= \operatorname{Re} \int_{b_1}^{a_1} (\xi_{1+}(s) - \xi_{2+}(s)) ds \\ &\quad + \operatorname{Re} \int_{a_1}^z (\xi_1(s) - \xi_2(s)) ds \\ &= - \int_z^{a_1} (\xi_1(s) - \xi_2(s)) ds > 0, \end{aligned} \quad (3.32)$$

where for the second equality we used (3.16) and for the final inequality we used (3.20).

Analogously,

$$2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 = 0, \quad z \in \Delta_2, \quad (3.33)$$

$$2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - \tilde{l}_3 = 0, \quad z \in \Delta_3, \quad (3.34)$$

$$2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - \tilde{l}_3 > 0, \quad z < \min\{a_*, a_1\}, \quad (3.35)$$

$$2U^{\mu_3}(a_*) + U^{\mu_1}(a_*) - U^{\mu_2}(a_*) = \tilde{l}_3. \quad (3.36)$$

If $\tau \geq \tau_c$, we evaluate (3.30), (3.33), (3.34) at the common point $a_* \in \Delta_1 \cap \Delta_2 \cap \Delta_3$ and take differences in order to get

$$l_1 - l_2 - \tilde{l}_3 = 0.$$

On the other hand, if $\tau < \tau_c$, we use (3.32) to get

$$2U^{\mu_1}(a_*) + U^{\mu_2}(a_*) + U^{\mu_3}(a_*) + \phi(a_*) > l_1.$$

We now combine this inequality with (3.33) evaluated at $z = a_* \in \Delta_2$ and (3.36), and conclude

$$\tilde{l}_3 > l_1 - l_2.$$

Hence, for both the precritical and supercritical cases we can define

$$l_3 := l_1 - l_2,$$

and with this definition, it follows that

$$2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - l_3 > 0, \quad z < \min\{a_*, a_1\},$$

$$2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - l_3 = 0, \quad z \in \Delta_3.$$

One more variational inequality is based on the fact that we deal with a critical trajectory of our quadratic differential. Recall that by (3.29),

$$2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 = \operatorname{Re} \psi_2(z), \quad z \in \mathbb{C} \setminus \Delta,$$

where

$$\psi_2(z) = \int_{a_2}^z (\xi_1(s) - \xi_3(s)) ds.$$

From (3.14), we know that ψ_2 is real-valued on $\Gamma \setminus \Delta_2$. Hence,

$$2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 = \psi_2(z), \quad z \in \Gamma \setminus \Delta_2.$$

Moreover, by (3.15), the real-valued function ψ_2 is monotone on each connected component of $\Gamma \setminus \Delta_2$. Analyzing at ∞ , we see that

$$\psi_2(z) = z^3 + \mathcal{O}(z),$$

hence ψ_2 tends to $+\infty$ along Γ . Since it is zero at the endpoints a_2, b_2 of the connected components of $\Gamma \setminus \Delta_2$, we conclude that ψ_2 is always positive on $\Gamma \setminus \Delta_2$, and hence

$$2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 > 0, \quad z \in \Gamma \setminus \Delta_2.$$

We summarize our findings in the following theorem:

Theorem 3.7. *Let measures μ_1, μ_2, μ_3 be defined in (3.21). Then there exist real constants l_1, l_2 and*

$$l_3 := l_1 - l_2,$$

such that the following variational conditions are satisfied:

$$\begin{aligned}
2U^{\mu_1}(z) + U^{\mu_2}(z) + U^{\mu_3}(z) + \phi(z) - l_1 &= 0, & z \in \Delta_1, \\
2U^{\mu_1}(z) + U^{\mu_2}(z) + U^{\mu_3}(z) + \phi(z) - l_1 &> 0, & z \in \Gamma_1 \setminus \Delta_1, \\
2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 &= 0, & z \in \Delta_2, \\
2U^{\mu_2}(z) + U^{\mu_1}(z) - U^{\mu_3}(z) + \phi(z) - l_2 &> 0, & z \in \Gamma_2 \setminus \Delta_2, \\
2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - l_3 &= 0, & z \in \Delta_3, \\
2U^{\mu_3}(z) + U^{\mu_1}(z) - U^{\mu_2}(z) - l_3 &> 0, & z \in \Gamma_3 \setminus \Delta_3.
\end{aligned} \tag{3.37}$$

where $\Gamma_1 = [a_*, +\infty)$, $\Gamma_2 = \Gamma$, $\Gamma_3 = (-\infty, a_*]$.

Remark 3.8. In the pre-critical case ($\tau < \tau_c$), $\Delta_3 = \emptyset$, thus the equality on Δ_3 in (3.37) is void.

The equalities in (3.37) are the same as those in (1.26). The extra information in Theorem 3.7 is coming from the remaining equations, which assure that the triplet (μ_1, μ_2, μ_3) is the (unique) minimizer of the energy functional (1.2) over measures (ν_1, ν_2, ν_3) satisfying (1.4) (with μ_j replaced by ν_j) and $\text{supp } \mu_j \subset \Gamma_j$, $j = 1, 2, 3$, see [13, Theorem 1.8]. This is equivalent to Theorem 1.16 with $\Gamma = \Gamma_2$.

4. GLOBAL STRUCTURE OF THE TRAJECTORIES IN THE CUBIC CASE

4.1. Dynamics of the singularities. A natural first step in the study of the structure of the trajectories of a quadratic differential is to clarify the position and the character of its singular points. In the case of the quadratic differential (3.12) we have to analyze the location and the dynamics of the branch points and the double point of the Riemann surface \mathcal{R} (corresponding to the spectral curve (3.9)) as functions of the parameter τ .

Recall that in Section 3.2 we have denoted the branch points of (3.9) by a_1, a_2, b_1 and b_2 , with the conventions

$$a_1, b_1 \in \mathbb{R}, \quad a_1 < b_1, \quad a_2 = \bar{b}_2, \quad \text{Im } a_2 < 0,$$

and the double point of (3.9) by b_* .

Proposition 4.1. *The main parameters of the Riemann surface \mathcal{R} corresponding to the algebraic equation (3.9) exhibit the following behavior: as τ grows from 0 to $1/4$,*

- (i) *the coefficient c in (3.4) increases monotonically from $-3^{5/3}/4$ to 0;*
- (ii)
 - *the branch point a_1 decreases monotonically from $3^{2/3}/4$ to $-\infty$;*
 - *the branch point b_1 increases monotonically from $3^{2/3}/4$ to $3^{2/3}/2$;*
 - *the double point b_* increases monotonically from $3^{-1/3}$ to $+\infty$.*
- (iii) *Always*

$$b_1 \leq b_*,$$

and the equality is attained only for $\tau = 1/12$.

Proof. Recall that we chose c in (3.4) to be real, thus (i) follows directly from the definition of c .

Any branch point $z = \zeta$ of \mathcal{R} is a zero of the polynomial q_1 of degree 4 and real coefficients, defined in (3.6), so that

$$\partial_\tau \zeta = -\frac{\partial_\tau q_1(\zeta)}{\partial_z q_1(\zeta)}. \tag{4.1}$$

It is easy to check that the resultant of q_1 and $\partial_\tau q_1$ (with respect to the variable z) is

$$\text{const} \times (1 - 4\tau)^{-\frac{5}{3}} (27 - 100\tau)^3 \neq 0 \quad \text{for } \tau \in (0, 1/4),$$

which implies that for any root ζ of q_1 , $\partial_\tau q_1(\zeta) \neq 0$, and thus preserves its sign in the whole range of values of $\tau \in (0, 1/4)$.

On the other hand, since a_1 (resp., b_1) is the smallest (resp., largest) real root of q_1 , and the leading coefficient of q_1 is positive, we know that

$$\partial_z q_1(a_1) < 0, \quad \partial_z q_1(b_1) > 0.$$

Straightforward calculations show that for $\tau = 1/8$, $a_1 = 0 < b_1$ and

$$\partial_\tau q_1(z) = -256u^4 - 64((u-1)^2 + 1), \quad u = \left(\frac{2}{3}\right)^{2/3} z.$$

In particular, $\partial_\tau q_1(z) \leq -64$, and we conclude that in this case,

$$\partial_\tau q_1(a_1) = -128, \quad \partial_\tau q_1(b_1) < 0.$$

Hence, $\partial_\tau q_1(a_1), \partial_\tau q_1(b_1) < 0$ for all $\tau \in (0, 1/4)$, and by (4.1), a_1 is a decreasing and b_1 is an increasing function of τ .

From the expression (3.7) it follows that

$$b_* = \frac{1}{3^{1/3}(1-4\tau)^{1/3}}.$$

Replacing it in (3.6) we get

$$q_1(b_*) = \frac{32(1-12t)^2}{9(1-4t)} > 0 \quad \text{for } \tau \in (0, 1/4), \quad \tau \neq 1/12.$$

Since for $\tau = 0$,

$$a_1 = b_1 = \frac{3^{2/3}}{4} < b_* = 3^{-1/3}, \quad (4.2)$$

this concludes the proof of (iii). Finally, it remains to observe that for $\tau = 1/4$,

$$q_1(z) = \frac{128}{9} \left(z^3 - \frac{9}{8} \right),$$

which has one real positive root ($3^{2/3}/2$) and two complex conjugate ones. This shows that

$$\lim_{\tau \rightarrow 1/4-} a_1 = -\infty, \quad \lim_{\tau \rightarrow 1/4-} b_1 = \frac{3^{2/3}}{2}.$$

□

We finish this section with a technical lemma that will be used later.

Proposition 4.2. *For $0 < \tau < 1/12$, the polynomial D in (3.2) does not have zeros on (a_1, b_1) , whereas for $1/12 < \tau < 1/4$, D has exactly one zero on (a_1, b_1) . Moreover, for $0 < \tau < 1/12$, $D(b_1) = 0$ only for $\tau = 1/12$, and $D(a_1)$ is never zero.*

Proof. On one hand, the discriminant of the polynomial D with respect to z is

$$f(\tau) = \text{const } \tau(1-12\tau)^2 (128\tau^2 - 32\tau - 1)$$

which shows that for $\tau \in (0, 1/4)$, polynomial D has no multiple roots as long as $\tau \neq 1/12$.

On the other hand, we have seen in Section 3.1 that the branch points of \mathcal{R} are the zeros of the polynomial q_1 defined in (3.6). The resultant (also w.r.t. z) of D and q_1 is

$$g(\tau) = \text{const}(1 - 12\tau)^3(27 - 100\tau)^3,$$

so again, for $\tau \in (0, 1/4)$, polynomials D and q_1 have no common roots as long as $\tau \neq 1/12$, in particular implying that D does not vanish at a_1, b_1 for $\tau \neq 1/12$. For $\tau = 1/12$, we compute $b_1 = b_* = 2^{-1/3}$ and factor

$$D(z) = -(z - b_*)^2(2z^4 + 2^{5/3}z^3 + 3 \cdot 2^{1/3}z^2 + z + 2^{-4/3}),$$

so for $\tau = 1/12$ we have $D(b_1) = 0$ and $D(a_1) < 0$.

It is worth pointing out that both f and g can be easily found by means of any computer algebra software.

Having in mind that D and q_1 do not share roots for $\tau \neq 1/12$, it is sufficient to establish the assertion concerning the zeros of D for a single value of τ in $(0, 1/12)$, and for a single value of τ in $(1/12, 1/4)$.

For $\tau = 1/4$, D is a quadratic polynomial in z^3 so its roots can be explicitly computed; we get that its only real roots are

$$z_1 = \left(\frac{3 - \sqrt{3}}{4}\right)^{1/3}, \quad z_2 = \left(\frac{3 + \sqrt{3}}{4}\right)^{1/3}.$$

Since

$$\left(\frac{3 - \sqrt{3}}{4}\right)^{1/3} < \frac{3^{2/3}}{2} < \left(\frac{3 + \sqrt{3}}{4}\right)^{1/3},$$

from Proposition 4.1 (ii) we see that for $\tau = 1/4 - \epsilon$, for a certain small value of $\epsilon > 0$,

$$a_1 < z_1 < b_1 < z_2.$$

On the other hand, it is easy to check that for $\tau = 0$,

$$D(a_1) = D(b_1) = D(3^{2/3}/4) < 0,$$

so that for $\tau = \epsilon$, for a certain small value of $\epsilon > 0$, D does not vanish on (a_1, b_1) . \square

4.2. The Riemann surface associated to the algebraic equation. In Section 3.2 we described the construction of the three-sheeted Riemann surface $\mathcal{R} = \mathcal{R}_1 \cup \mathcal{R}_2 \cup \mathcal{R}_3$ and of the branch cuts in such a way that the three solutions ξ_j of (3.9), specified by the asymptotic conditions (3.10), become meromorphic on the respective sheet \mathcal{R}_j , with poles only at $z = \infty$. They are also pairwise distinct as long as $\text{Discr}(z)$, defined in (3.5), does not vanish, i.e. for $z \notin \{a_1, b_1, a_2, b_2, b_*\}$.

The arc Δ_2 intersects \mathbb{R} in a unique point a_* ; there is a critical value $\tau_c \in (1/12, 1/4)$, to be specified later, such that $a_* < a_1$ for $0 < \tau < \tau_c$ (what we called the precritical regime), and $a_* \in (a_1, b_1)$ for $\tau_c < \tau < 1/4$ (the supercritical regime).

As a first result we establish some relations between the solutions $\xi_j(z)$.

Proposition 4.3. *Let $\tau \in (0, 1/4)$, $\tau \neq \tau_c$. Then*

(i) for $x \in \mathbb{R} \setminus (\Delta_1 \cup \Delta_2 \cup \Delta_3)$,

$$\xi_3(x) < \xi_2(x) < \xi_1(x), \quad x < \min(a_*, a_1), \quad (4.3)$$

$$\xi_2(x) < \xi_3(x) < \xi_1(x), \quad x > b_*, \quad (4.4)$$

$$\xi_1(x) < \xi_2(x) < \xi_3(x), \quad a_* < x < a_1, \quad \text{if } \tau < \tau_c, \quad (4.5)$$

$$\xi_2(x) < \xi_1(x) < \xi_3(x), \quad b_1 < x < b_*, \quad \text{if } 0 < \tau < 1/12, \quad (4.6)$$

$$\xi_3(x) < \xi_2(x) < \xi_1(x), \quad b_1 < x < b_*, \quad \text{if } 1/12 < \tau < 1/4. \quad (4.7)$$

Additionally,

$$\begin{aligned} \xi_2(b_*) &< \xi_3(b_*) = \xi_1(b_*), \quad \text{for } 0 < \tau < 1/12, \\ \xi_2(b_*) &= \xi_3(b_*) < \xi_1(b_*), \quad \text{for } 1/12 < \tau < 1/4. \end{aligned} \quad (4.8)$$

(ii) on $\Delta_1 \cup \Delta_2 \cup \Delta_3$,

- for $x \in \overset{\circ}{\Delta}_1 := \Delta_1 \setminus \{\max(a_1, a_*), b_1\}$,

$$\xi_2(x) = \overline{\xi_1(x)} \in \mathbb{C} \setminus \mathbb{R}, \quad \xi_3(x) \in \mathbb{R}, \quad (4.9)$$

and

$$\xi_{1\pm}(x) = \xi_{2\mp}(x), \quad \xi_{3+}(x) = \xi_{3-}(x). \quad (4.10)$$

- for $z \in \overset{\circ}{\Delta}_2 := \Delta_2 \setminus \{a_2, b_2, \max(a_1, a_*)\}$,

$$\xi_{1\pm}(z) = \xi_{3\mp}(z), \quad \xi_{2+}(z) = \xi_{2-}(z). \quad (4.11)$$

- for $x \in \overset{\circ}{\Delta}_3 := \Delta_3 \setminus \{a_1, a_*\}$ (when $\tau > \tau_c$),

$$\xi_2(x) = \overline{\xi_3(x)} \in \mathbb{C} \setminus \mathbb{R}, \quad \xi_1(x) \in \mathbb{R}, \quad (4.12)$$

and

$$\xi_{2\pm}(z) = \xi_{3\mp}(z), \quad \xi_{1+}(z) = \xi_{1-}(z). \quad (4.13)$$

Moreover,

$$\xi_1(a_1) = \xi_2(a_1), \quad \text{if } \tau < \tau_c, \quad (4.14)$$

$$\xi_3(a_1) = \xi_2(a_1), \quad \text{if } \tau > \tau_c, \quad (4.15)$$

$$\xi_1(b_1) = \xi_2(b_1), \quad (4.16)$$

$$\xi_1(a_2) = \xi_3(a_2) \quad \text{and} \quad \xi_1(b_2) = \xi_3(b_2). \quad (4.17)$$

Proof. Recall that we deal with the case $\alpha \in (0, 1/2)$, so that $0 < \alpha < 1 - \alpha < 1$. The behavior at infinity in (3.10) gives (4.3) and also (4.4). Furthermore, we have established in Proposition 4.1 that if $\tau \neq 1/12$ then $b_* \notin \{a_1, b_1, a_2, b_2\}$, and in this case b_* is a double zero of $\text{Discr}(z)$. This means that when $\tau \neq 1/12$, only two of the three values $\xi_1(b_*)$, $\xi_2(b_*)$, $\xi_3(b_*)$ coincide, and since

$$\xi_1(b_*) + \xi_2(b_*) + \xi_3(b_*) = 0,$$

the coincident two differ in sign from the third one. Furthermore, since also

$$(\xi_1 \xi_2 \xi_3)(b_*) = -D(b_*),$$

we see that the sign of this third one is opposite to the sign of $D(b_*)$. But direct calculations show that

$$D(b_*) = -\frac{(12\tau - 1)^3}{36(1 - 4\tau)^2} \begin{cases} > 0, & \text{for } 0 < \tau < 1/12, \\ < 0, & \text{for } 1/12 < \tau < 1/4, \end{cases}$$

which means that for $0 < \tau < 1/12$ the two coincident values of $\xi_j(b_*)$ are the largest two, and the other way around if $1/12 < \tau < 1/4$. It remains to use (4.4) in order to establish (4.8).

Inequalities (4.6)–(4.7) follow by noticing that b_* is a double zero of $\text{Discr}(z)$, so it is a simple zero of $\xi_1 - \xi_3$ (for $0 < \tau < 1/12$) and of $\xi_2 - \xi_3$ (for $1/12 < \tau < 1/4$).

From Proposition 4.2 we know that $D(a_1) \neq 0$, $D(b_1) \neq 0$ for $\tau \neq 1/12$. In particular, since for $\tau = 0$ (see (4.2)), $a_1 = b_1 = 3^{2/3}/4$ and $D(a_1) = D(b_1) = -81/2048 < 0$, we conclude that $D(a_1) < 0$ and $D(b_1) < 0$ for $0 \leq \tau < 1/12$. In other words, the smallest two of the three solutions ξ_j come together at a_1 and at b_1 . Now (4.6) implies (4.16) for $0 \leq \tau < 1/12$.

On the other hand, recall that by Proposition 4.1, $D(a_1) \rightarrow -\infty$ as $\tau \rightarrow 1/4$, and that for $\tau = 1/4$,

$$D(z) = -2z^6 + 3z^3 - \frac{3}{4}, \quad b_1 = 3^{2/3}/2,$$

so that $D(b_1) = 3/32 > 0$. Using the same arguments we conclude that for $1/12 < \tau < 1/4$, the largest two of the three solutions ξ_j come together at b_1 (and (4.7) yields (4.16)), and the smallest two become confluent at a_1 (and (4.15) follows from (4.3)).

The discriminant $\text{Discr}(x)$ is negative for $a_1 < x < b_1$, so just one of the solutions ξ_j is real on $a_1 < x < b_1$, and the other two are complex-conjugates of each other. Since we have ruled out the coincidence of the three branches, equality (4.16) implies that ξ_1, ξ_2 are non-real on Δ_1 , which yields (4.9)–(4.10). By the same argument we also have (4.14), as well as (now using (4.3)) the identities (4.12)–(4.13).

Recall that we already established that the smallest two of the three solutions ξ_j come together at a_1 , and (4.14) shows that for $\tau < \tau_c$,

$$\xi_1(x) < \xi_3(x), \quad \xi_2(x) < \xi_3(x), \quad \text{for } a_* < x < a_1. \quad (4.18)$$

For $\tau < \tau_c$, when crossing Δ_2 two of the three solutions ξ_j are swapped, and the third one remains invariant. But the only option compatible both with (4.3) and (4.18) is (4.5).

Finally, for $\tau < \tau_c$, the inequalities (4.3) and (4.5) show that ξ_2 is continuous across Δ_2 , and as a consequence (4.11) has to hold, obviously for the full range of $\tau \in (0, 1/4)$. This, in turn, implies the last identity (4.17). \square

Proposition 4.1 gives the formal proof of the fact that the Riemann surface \mathcal{R} , described in Section 3.2, is actually the Riemann surface of the cubic equation (3.9), and with this construction the function ξ_j is meromorphic on \mathcal{R}_j and satisfies the asymptotic expansion (3.10). Moreover, $\xi : \mathcal{R} \rightarrow \mathbb{C}$, given by $\xi \equiv \xi_j$ on \mathcal{R}_j , is meromorphic on \mathcal{R} , giving the global solution to the algebraic equation (3.9). From Theorem 1.12, or alternatively Proposition 3.1, the Riemann surface \mathcal{R} has genus 0.

Let $\pi : \mathcal{R} \rightarrow \overline{\mathbb{C}}$ be the canonical projection on the Riemann surface \mathcal{R} . For a point $p \in \overline{\mathbb{C}} \setminus (\Delta_1 \cup \Delta_2 \cup \Delta_3)$, we denote by $p^{(j)}$, $j = 1, 2, 3$, its preimage by π on \mathcal{R}_j , that is,

$$\{p^{(j)}\} = \pi^{-1}(\{p\}) \cap \mathcal{R}_j, \quad j = 1, 2, 3.$$

This notation is trivially extended to $p \in \Delta_1 \cup \Delta_2 \cup \Delta_3$ by taking boundary values.

Notice that for the branch points it is valid

$$b_1^{(1)} = b_1^{(2)}, \quad a_2^{(2)} = a_2^{(3)}, \quad b_2^{(2)} = b_2^{(3)},$$

and

$$a_1^{(1)} = a_1^{(2)} \quad (\tau \leq \tau_c), \quad a_1^{(2)} = a_1^{(3)} \quad (\tau > \tau_c),$$

whereas if p belongs to a cut connecting exactly two sheets, say \mathcal{R}_j and \mathcal{R}_k , then

$$p_{\pm}^{(j)} = p_{\mp}^{(k)}.$$

We insist that the construction of \mathcal{R} is independent of the concrete choice of the cut Δ_2 ; this freedom will be used latter to specify an appropriate Δ_2 . Namely, in the next section we will show that Δ_2 can be made coincident with a critical trajectory of the quadratic differential ϖ , defined in (3.12), which connects a_2 and b_2 , see Definition 4.5 below.

4.3. Computation of width parameters. Certain integrals of the function Q defined in (3.11) will play a crucial role in the upcoming analysis of the dynamics of the trajectories of the quadratic differential $\varpi = -Q^2(z) dz^2$ on \mathcal{R} . They can be also formulated in terms of certain Abelian integrals on \mathcal{R} .

Namely, we are interested in

$$\begin{aligned} \omega_1 &= \omega_1(\tau) = \operatorname{Re} \int_{b_2^{(1)}}^{a_1^{(1)}} Q(s) ds = \operatorname{Re} \int_{b_2}^{a_1} (\xi_2(s) - \xi_3(s)) ds, \\ \omega_2 &= \omega_2(\tau) = \operatorname{Re} \int_{b_2^{(2)}}^{a_1^{(2)}} Q(s) ds = \operatorname{Re} \int_{b_2}^{a_1} (\xi_1(s) - \xi_3(s)) ds, \\ \omega_3 &= \omega_3(\tau) = \operatorname{Re} \int_{b_2^{(3)}}^{a_1^{(3)}} Q(s) ds = \operatorname{Re} \int_{b_2}^{a_1} (\xi_1(s) - \xi_2(s)) ds, \\ \omega_4 &= \omega_4(\tau) = \operatorname{Re} \int_{b_2^{(1)}}^{b_*^{(1)}} Q(s) ds = \operatorname{Re} \int_{b_2}^{b_*} (\xi_2(s) - \xi_3(s)) ds, \end{aligned} \tag{4.19}$$

with $\omega_4(\tau)$ defined for $\tau \geq 1/12$. In this definition we understand that we integrate between two points on a sheet \mathcal{R}_j along a path that stays entirely in \mathcal{R}_j .

The values ω_j are correctly defined regardless of the precise choice of the integration paths. Indeed, the residues of the functions ξ_1, ξ_2, ξ_3 at ∞ are real (and independent of the value of τ), see (3.10), so the integral of Q along a big loop on either \mathcal{R}_j encircling $\infty^{(j)}$ is purely imaginary. It remains to notice that the genus of \mathcal{R} is zero, so that any closed contour around either branch cut Δ_j can be deformed to such a big loop.

Analytic computation of ω_j 's is a formidable task. Instead, we have computed them numerically, see Figure 4 for the result. Since the integrands in (4.19) are multivalued functions, the numerical integration requires to implement an analytic continuation of each branch of ξ . We give further details in the Appendix C.

The functions ω_1, ω_2 , have two and one zeros on $(0, 1/4)$, respectively, whereas ω_3 and ω_4 do not vanish on the intervals $(0, 1/4)$ and $(1/12, 1/4)$, respectively.

We denote the zeros of ω_1 by τ_1, τ_2 and of ω_2 by τ_c . We have

$$\tau_1 \approx 0.12487351, \quad \tau_c \approx 0.1913565, \quad \tau_2 \approx 0.2289555,$$

so that they satisfy

$$0 < \frac{1}{12} < \tau_1 < \tau_c < \tau_2 < 1/4. \tag{4.20}$$

We point out that the value τ_c is the one used in Section 3.2 in the construction of the Riemann surface \mathcal{R} , and it is equivalently determined by (1.34).

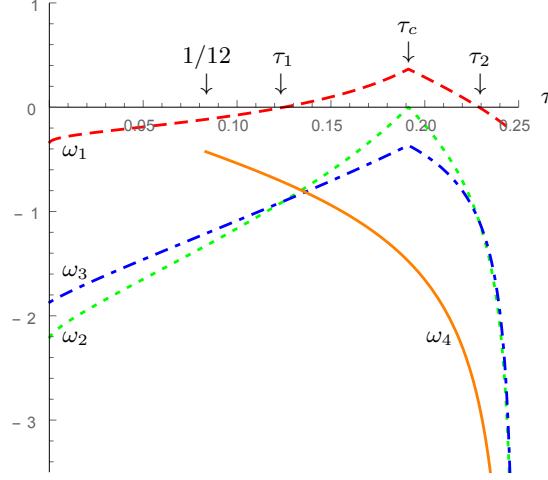


FIGURE 4. The graphics of the functions ω_1 (upper dashed line), ω_2 (short dashed line), ω_3 (dashed line with dots) and ω_4 (continuous line).

It is also clear from Figure 4 that $\omega_1 \neq \omega_4$ for $\tau > 1/12$.

4.4. Critical points of the quadratic differential. We follow the construction (1.28), (3.12) and define the meromorphic quadratic differential $\varpi = -Q^2(z) dz^2$ on the Riemann surface \mathcal{R} .

We start by analyzing the character of the critical points of ϖ . For instance, the local parameter at $z = a_1^{(1)}$ is

$$z = a_1^{(1)} + u^2,$$

so that $dz^2 = 4u^2 du^2$. In consequence, $(\xi_j - \xi_k)^2(z) dz^2$ has a double zero at $z = a_1^{(1)}$ if and only if $\xi_j(a_1^{(1)}) \neq \xi_k(a_1^{(1)})$. A similar analysis at the rest of the points of \mathcal{R} yields the following classification of the critical points of ϖ on the Riemann surface (see Figure 5):

- (a) For $0 < \tau < \tau_c$:
 - (i) Simple zeros at $a_1^{(3)}, b_1^{(3)}, a_2^{(2)}, b_2^{(2)}$;
 - (ii) Double zeros at $a_1^{(1)}, b_1^{(1)}$ (only for $\tau \neq 1/12$), $a_2^{(1)}, b_2^{(1)}$;
 - (iii) Double zero at $b_*^{(2)}$ if $\tau < 1/12$ and a double zero at $b_*^{(1)}$ if $\tau > 1/12$;
 - (iv) Zero of order 4 at $b_1^{(1)} = b_*^{(1)}$ if $\tau = 1/12$;
 - (v) Double pole at $\infty^{(1)}$ with a real residue;
 - (vi) Poles of order 8 at $\infty^{(2)}, \infty^{(3)}$.
- (b) For $\tau_c \leq \tau < 1/4$:
 - (i) Simple zeros at $a_1^{(1)}, b_1^{(3)}, a_2^{(2)}, b_2^{(2)}$;
 - (ii) Double zeros at $a_1^{(2)}, b_1^{(1)}, a_2^{(1)}, b_2^{(1)}$;
 - (iii) Double zero at $b_*^{(1)}$;
 - (iv) Double pole at $\infty^{(1)}$ with a real residue;
 - (v) Poles of order 8 at $\infty^{(2)}, \infty^{(3)}$.

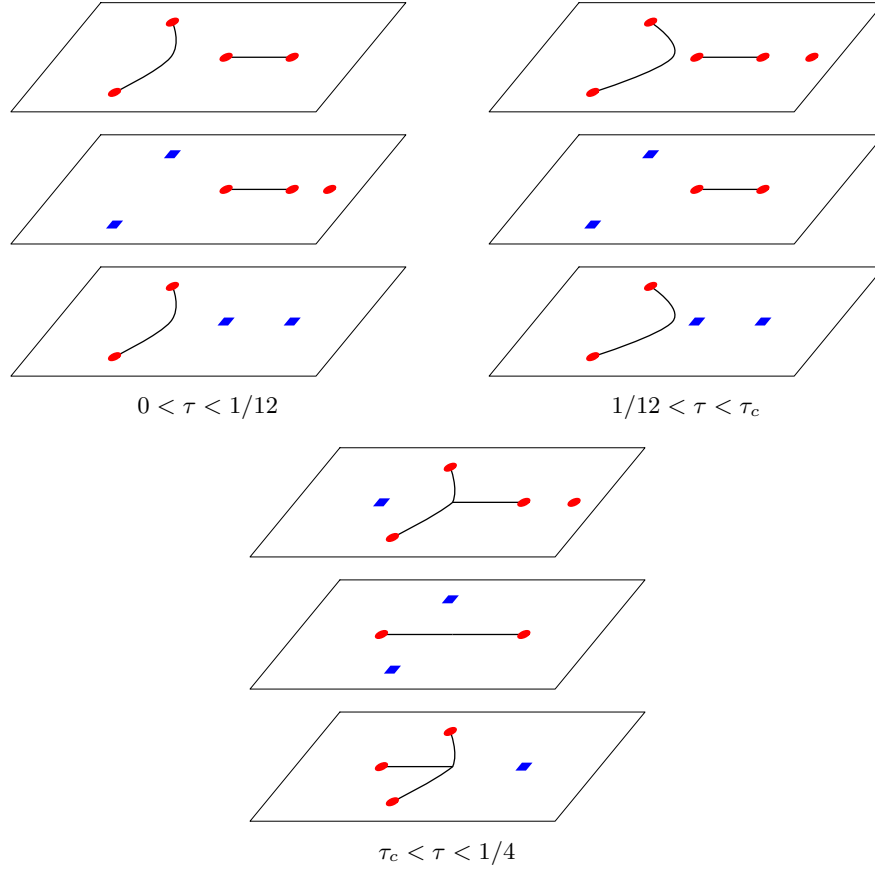


FIGURE 5. The finite critical points of ϖ in each case. Simple and double zeros are represented, respectively, with squares and dots.

It is instructive to think of the critical points as evolving dynamically with τ . Under this perspective, Propositions 4.1 and 4.3 show that for τ small the double point (node) b_* corresponds to the double zero $b_*^{(2)}$ on the second sheet. When $\tau = 1/12$, the points b_* , b_1 coalesce, giving rise to a higher order zero at the branch point $b_1^{(1)} = b_1^{(2)}$. For larger values of τ , the double point emerges on the first sheet: now $b_*^{(2)}$ is a regular point whereas $b_*^{(1)}$ is a double zero.

In the same spirit, the simple zero a_1 of the discriminant of (3.9) carries two critical points of ϖ . For $\tau < \tau_c$ these are the simple zero $a_1^{(3)}$ and the double zero $a_1^{(1)} = a_1^{(2)}$. For values of τ larger than τ_c , these points interchange their roles: $a_1^{(1)}$ is a simple zero and $a_1^{(3)} = a_1^{(2)}$ is a double zero.

4.5. Analyzing the global structure of trajectories.

4.5.1. *General principles.* The rest of this section is devoted to the description of the critical graph of the quadratic differential (3.12) for the whole range $0 \leq \tau < 1/4$. One of the outcomes of our analysis is the following theorem:

Theorem 4.4. *For the quadratic differential $\varpi = -Q^2(z) dz^2$ and for all values of the parameter $0 \leq \tau \leq \tau_c$ there exists a critical trajectory of ϖ joining $a_2^{(2)}$ and $b_2^{(2)}$ on \mathcal{R}_2 , whose projection by π on \mathbb{C} is a real-symmetric analytic arc Δ_2 joining a_2 and b_2 .*

For $\tau_c \leq \tau \leq 1/4$, there is an arc of critical trajectory of ϖ joining $a_2^{(2)}$ with a point $a_^{(2)}$ on the interval $[a_1^{(2)}, b_1^{(2)}]$ which is determined by (1.35), and the conjugate symmetric arc of trajectory joining $b_2^{(2)}$ with the same point. The projection by π on \mathbb{C} of the union of these two arcs of trajectories is also denoted by Δ_2 .*

Recall that τ_c was formally introduced in Section 4.3. In virtue of the results in Section 3.2, in particular Corollary 3.6, Theorem 4.4 implies Theorem 1.14.

We remind the reader that up to now the branch cut, separating the sheets \mathcal{R}_1 and \mathcal{R}_3 , was free (see Figure 2). In what follows we agree in the following:

Definition 4.5. The curve, connecting a_2 and b_2 as part of the branch cut separating the sheets \mathcal{R}_1 and \mathcal{R}_3 , is always given by the lift of Δ_2 (from Theorem 4.4) to the sheets \mathcal{R}_1 and \mathcal{R}_3 .

In the next sections we will show that this definition is consistent with our construction of the Riemann surface \mathcal{R} .

One important fact is that the residues of ϖ at the poles at infinity (and the local behavior of the trajectories there) are independent of τ : at $\infty^{(1)}$ they are closed analytic curves (so that $\infty^{(1)}$ is the center of a circle domain, see Appendix B), while $\infty^{(2)}, \infty^{(3)}$ attract trajectories in 6 asymptotic directions, given by the angles

$$\theta_j^{(\infty)} = \frac{2j-1}{6} \pi, \quad j = 1, \dots, 6. \quad (4.21)$$

Critical values (4.20) split the interval $(0, 1/4)$ into the subintervals $(0, 1/12)$, $(1/12, \tau_1)$, (τ_1, τ_c) , (τ_c, τ_2) and $(\tau_2, 1/4)$. We will show that the topology of the critical graph remains invariant in each of these intervals.

The methodology we use can be summarized as follows:

- (i) Compute the critical graph for τ equal to one of the critical values (4.20).
- (ii) Analyze the possible deformation of the trajectories for the values $\tau + \varepsilon$, with $\varepsilon > 0$ small, identifying the trajectories that display a phase transition.
- (iii) Prove that the topology of the critical graph is invariant inside the subinterval of interest, by analyzing the behavior of the widths ω_j 's and showing that the corresponding strip and ring domains can not disappear.

Along the way, will use some general *principles* that we enumerate here:

- P.1** Quadratic differential ϖ has no recurrent trajectory for any value of τ , see Jenkin's Three Poles Theorem in Section B.3.
- P.2** If γ is an arc of trajectory of ϖ , then $\bar{\gamma}$, corresponding to the lift of the complex conjugate of $\pi(\gamma)$ to the same sheet, is also an arc of trajectory.
- P.3** The complement of the critical graph of ϖ on \mathcal{R} cannot have a simply connected component without poles on its boundary: that would contradict Corollary B.3 or the maximum principle for harmonic functions on a compact Riemann surface.

- P.4** The meromorphic function Q^2 depends analytically on the parameter τ . Hence the critical graph of ϖ , and in particular all its critical trajectories, depend continuously (in any reasonable topology, for instance, in the Hausdorff distance) from τ .
- P.5** If for a certain value $\tau = A$, the point p belongs to the half plane domain for $\infty^{(k)}$ determined by the angles $\theta_j^{(\infty)}, \theta_{j+1}^{(\infty)}$, then the same holds true for a small neighborhood of values $\tau \in (A - \varepsilon, A + \varepsilon)$, $\varepsilon > 0$. The point p is also allowed to depend continuously on τ .
- P.6** If for a certain value $\tau = A$, an arc of trajectory emanating from a given point p intersects the real line at a *regular* point, then the same holds true for $\tau \in (A - \varepsilon, A + \varepsilon)$, $\varepsilon > 0$. As before, the point p is allowed to depend continuously on τ .

There will be one more useful tool that we will employ several times in our analysis, formulated as Proposition 4.7 below.

When describing the structure and the evolution of the trajectories of the quadratic differential ϖ we face the dilemma of either a formalization of each statement, with a precise formulation of the behavior of every trajectory in every situation, or a much more visual description, with rigorous proofs but illustrated by a number of figures. We opted for the second choice².

Next, we agree on some convention about trajectories. Let $p^{(j)} \in \mathcal{R}_j$ be a zero of ϖ . We denote by $\gamma_1(p^{(j)}), \gamma_2(p^{(j)}), \dots$, the trajectories of ϖ *emanating* from p on \mathcal{R}_j , in such a way that their canonical projections $\pi(\gamma_n(p^{(j)}))$, see Section 4.2, are enumerated in an anti-clockwise direction starting from the positive OX semiaxis³. Notice that when $p^{(j)}$ is a branch point of \mathcal{R} , so that $p^{(j)} = p^{(k)}$ for some $j \neq k$, trajectories $\gamma_n(p^{(j)})$ and $\gamma_n(p^{(k)})$ are different because they emerge from $p^{(j)} = p^{(k)}$ on different sheets of \mathcal{R} . Otherwise, when $p^{(j)}$ belongs to a single sheet \mathcal{R}_j , we occasionally drop the superindex (j) when it cannot lead us into confusion.

Given two points $p, q \in \mathbb{R}$, the integral

$$\int_p^q \sqrt{-\varpi}$$

along a contour γ connecting p and q is understood to be the integral of any analytic continuation of the meromorphic differential $\sqrt{-\varpi}$ along γ . This is well defined up to the branch of the square root, which will be clear in each context.

4.5.2. *Degenerate case* $\tau = 0$. For $\tau = 0$, the algebraic equation (3.9) reduces to

$$(\xi + z^2) \left(4\xi^2 - 8z^4 - 4\xi z^2 + 12z - 3 \times 3^{2/3} \right) = 0,$$

² We confess we might have been influenced by the famous quote of Vladimir Arnold [82]:

It is almost impossible for me to read contemporary mathematicians who, instead of saying “Petya washed his hands,” write simply: “There is a $t_1 < 0$ such that the image of t_1 under the natural mapping $t_1 \mapsto \text{Petya}(t_1)$ belongs to the set of dirty hands, and a $t_2, t_1 < t_2 \leq 0$, such that the image of t_2 under the above-mentioned mapping belongs to the complement of the set defined in the preceding sentence.”

³ This notation is not correctly defined only if the direction of a trajectory coincides with a branch cut, situation that will be explicitly avoided in what follows.

whose solutions, denoted in accordance with (3.10), are

$$\xi_1(z) = \frac{z^2 + \sqrt{3}\sqrt{h(z)}}{2}, \quad \xi_2(z) = -z^2, \quad \xi_3(z) = \frac{z^2 - \sqrt{3}\sqrt{h(z)}}{2},$$

where the branch of the square root is chosen to be positive for large real values,

$$h(z) = 3z^4 - 4z + 3^{2/3} = 3(z - a_2)(z - b_2)(z - b_*)^2,$$

and the points a_1, b_1, a_2, b_2, b_* are given explicitly by

$$a_1 = b_1 = \frac{3^{2/3}}{4}, \quad b_2 = \overline{a_2} = \frac{1}{3^{1/3}}(-1 + i\sqrt{2}), \quad b_* = \frac{1}{3^{1/3}}.$$

The cut Δ_1 is reduced to a single point a_1 , and the sheet \mathcal{R}_2 is detached from the others. Since (3.9) is reducible, its Riemann surface is in fact the union of two Riemann surfaces,

$$\mathcal{R}_2 = \overline{\mathbb{C}} \quad \text{and} \quad \tilde{\mathcal{R}} = \mathcal{R}_1 \cup \mathcal{R}_3.$$

Here $\mathcal{R}_1 = \mathcal{R}_3 = \overline{\mathbb{C}} \setminus \Delta_2$, and Δ_2 is a simple curve connecting the points a_2, b_2 , to be precisely specified later. The quadratic differential (3.11) degenerates into two quadratic differentials ϖ_1 on $\overline{\mathbb{C}}$ and ϖ_2 on $\tilde{\mathcal{R}}$, namely

$$\begin{aligned} \varpi_1 &= -3h(z)dz^2 \quad \text{on } \overline{\mathbb{C}}, \\ \varpi_2 &= \begin{cases} -\frac{1}{4}(3z^2 - \sqrt{3}\sqrt{h(z)})^2 dz^2 & \text{on } \mathcal{R}_1, \\ -\frac{1}{4}(3z^2 + \sqrt{3}\sqrt{h(z)})^2 dz^2 & \text{on } \mathcal{R}_3. \end{cases} \end{aligned} \quad (4.22)$$

We analyze the structure of their critical graphs next.

Trajectories of ϖ_1 , whose only critical points are as follows:

- Simple zeros at $z = a_2, b_2$;
- Double zero at $z = b_*$;
- Pole of order 8 at $z = \infty$.

Under the change of variables $z \mapsto \frac{i}{\sqrt[3]{3}}z$, the quadratic differential ϖ_1 becomes $-q(z)dz^2$, where

$$q(z) = -(z^2 - 2iz - 3)(z + i)^2$$

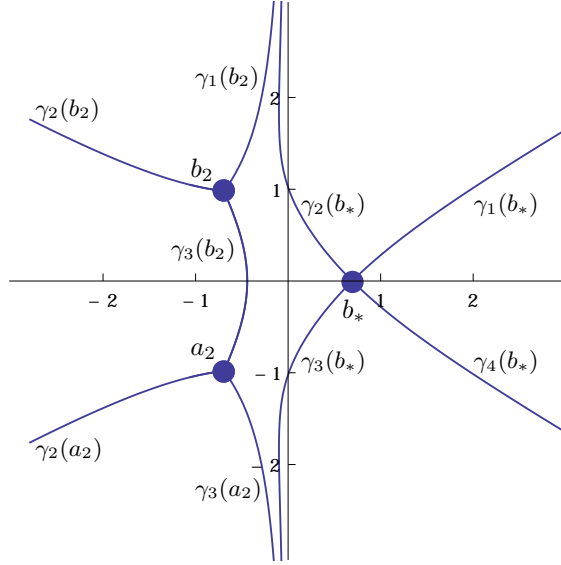
is, up to a multiplicative factor $\frac{1}{4}$, the same polynomial obtained in [27, eq. (2.1)]. Having in mind this identification, it was proven in [27, Theorem 2.1] that the trajectory $\gamma_3(b_2)$ of ϖ_1 connects b_2 and a_2 : in the notation introduced above, $\gamma_3(b_2) = \gamma_1(a_2)$, see Figure 6.

From [27] we also know that $\gamma_3(b_2)$ intersects the real axis at a point $a_* < b_*$, which can be calculated numerically: by

$$a_* \approx -0.441782. \quad (4.23)$$

The rest of the critical graph of ϖ_1 is as follows. Notice that due to the symmetry, we only need to describe the trajectories in the upper half plane.

The trajectories $\gamma_j(b_*)$, $\gamma_j(b_2)$, $j = 1, 2$, cannot be finite, see Principle **P.3** above. Hence, they all diverge to ∞ along the asymptotic directions (4.21), and according Theorem B.1, all directions are represented. There are 3 asymptotic directions for 4 trajectories in the upper half plane, so necessarily the divergence angle for $\gamma_1(b_*)$ is $\theta_1^{(\infty)}$, while the divergence angle for $\gamma_2(b_2)$ is $\theta_3^{(\infty)}$. Since two consecutive trajectories emanating from a zero cannot diverge to ∞ in the same direction (this would contradict Theorem B.2), we conclude that both $\gamma_2(b_*)$ and $\gamma_1(b_2)$ must diverge in the direction $\theta_2^{(\infty)}$, see Figure 6.

FIGURE 6. $\tau = 0$: the critical graph of ϖ_1 .

Trajectories of ϖ_2 , whose only critical points are as follows:

- Double zeros at $z = a_1^{(3)}$, $z = b_2^{(1)}$ and $z = a_2^{(1)}$;
- Double pole at $z = \infty^{(1)}$ with real residue;
- Pole of order 8 at $\infty^{(3)}$.

The double zeros $b_2^{(1)} = b_2^{(3)}$, $a_2^{(1)} = a_2^{(3)}$ are also branch points of $\tilde{\mathcal{R}}$, and the critical graph of ϖ_2 is made of trajectories $\gamma_j(a_2^{(1)})$, $\gamma_j(b_2^{(1)})$, $j = 1, 2$, emanating on \mathcal{R}_1 , and of trajectories $\gamma_j(a_2^{(3)})$, $\gamma_j(b_2^{(3)})$, $j = 1, 2$, along with $\gamma_j(a_1^{(3)})$, $j = 1, \dots, 4$, emanating on \mathcal{R}_3 , see Figure 7.

The branch cut Δ_2 connecting the sheets \mathcal{R}_1 and \mathcal{R}_3 , so far arbitrary, is chosen as $\Delta_2 = \gamma_3(b_2)$, where $\gamma_3(b_2)$ is the critical trajectory of ϖ_1 connecting b_2 to a_2 , as described above.

Lemma 4.6. *With the branch cut Δ_2 specified above, the critical trajectories $\gamma_j(b_2^{(1)})$ and $\gamma_j(a_2^{(1)})$ of ϖ_2 belong entirely to the sheet \mathcal{R}_1 .*

Proof. Suppose that one of the trajectories $\gamma_j(b_2^{(1)})$ emanating from b_2 intersects the cut Δ_2 for the first time at a point x . Clearly, this point (actually, its canonical projection) must lie in the upper half plane: otherwise we readily get that $\gamma_j(b_2^{(1)}) = \gamma_j(a_2^{(1)})$ and no intersection with Δ_2 occur.

Integrating from b_2 to x along $\gamma_j(b_2^{(1)})$ and using the definition of a trajectory we get

$$0 = \operatorname{Re} \int_{b_2^{(1)}}^{x^{(1)}} \sqrt{-\varpi} = \frac{3}{4} \operatorname{Re} \int_{b_2}^x s^2 ds - \frac{\sqrt{3}}{4} \operatorname{Re} \int_{b_2}^x \sqrt{h(s)}_{\pm} ds \quad (4.24)$$

where the \pm sign in the last integrand depends on the side of the cut Δ_2 to which x belongs. However, Δ_2 projects onto the trajectory $\gamma_3(b_2)$ of ϖ_1 , so the second

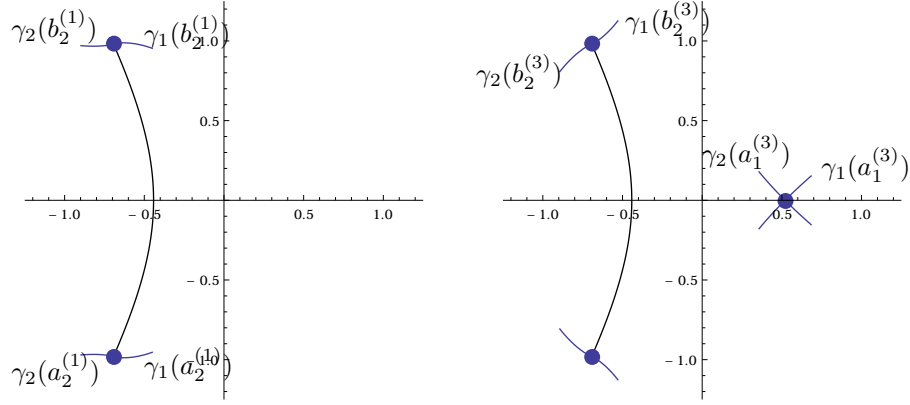


FIGURE 7. On the left it is represented the critical trajectories of ϖ_2 on \mathcal{R}_1 and on the right its critical trajectories on \mathcal{R}_3 , all of them in blue and just locally at the critical points. The black curve is the cut Δ_2 connecting \mathcal{R}_1 and \mathcal{R}_3 .

integral in the right-hand side of (4.24) is purely imaginary. Hence, this equation reduces to

$$\operatorname{Re} x^3 = \operatorname{Re} b_2^3 = \frac{5}{3}. \quad (4.25)$$

It was proved in [27] that the trajectory Δ_2 is contained in the domain bounded by the triangle with vertices b_2 , $3^{-1/3}(\sqrt{2} - 1)$ and a_2 . In particular, the part of Δ_2 on the upper half plane, and hence x , is contained in the domain bounded by the triangle T determined by the vertices b_2 , $\operatorname{Re} b_2$ and $3^{-1/3}(\sqrt{2} - 1)$. The function $z \mapsto \operatorname{Re} z^3$ has a unique maximum on T at the point $z = b_2$. Since $\operatorname{Re} z^3$ is harmonic, it cannot attain a maximum on the domain bounded by T , hence (4.25) can only occur if $x = b_2$, showing that $\gamma_j(b_2^{(1)}) \setminus \{b_2^{(1)}\}$ does not intersect the cut Δ_2 . \square

Recall that $\infty^{(1)}$ is the center of a circle domain, which means that all trajectories of ϖ_2 passing through sufficiently distant points on \mathcal{R}_1 are closed Jordan curves. In particular, no trajectory diverges to $\infty^{(1)}$, and every trajectory entirely contained in \mathcal{R}_1 has to be closed. Consequently, both trajectories $\gamma_j(b_2^{(1)})$ are closed as well, $\gamma_j(b_2^{(1)}) = \gamma_j(a_2^{(1)})$, and $\gamma_1(b_2^{(1)})$ (respectively $\gamma_2(b_2^{(1)})$) intersects the real line, say at the point c_* (respectively d_*). We claim that

$$d_* < a_* < c_* < a_1, \quad (4.26)$$

with a_* defined in (4.23).

Indeed, both c_* and d_* cannot lie on the same side of the cut Δ_2 without running into contradiction with the general principle **P.3** above. Hence, either $d_* < a_* < c_*$ or $c_* < a_* < d_*$. The latter is impossible without $\gamma_1(b_2^{(1)})$ and $\gamma_2(b_2^{(1)})$ intersecting somewhere in the upper half plane, which again contradicts **P.3**. We conclude that $d_* < a_* < c_*$.

Let us prove the inequality $c_* < a_1$. Using the definition of trajectory and (4.19) we get

$$0 = \operatorname{Re} \int_{b_2^{(1)}}^{c_*^{(1)}} Q(s) ds = \operatorname{Re} \int_{b_2}^{c_*} (\xi_2(s) - \xi_3(s)) ds = \omega_3 + \int_{a_1}^{c_*} (\xi_2(s) - \xi_3(s)) ds.$$

But $\omega_3 = \omega_3(0) < 0$, see Figure 4, so that

$$\int_{a_1}^{c_*} (\xi_2(s) - \xi_3(s)) ds > 0.$$

Function $\xi_2 - \xi_3$ is continuous and non-vanishing on $(a_*, +\infty)$, and by (3.10) it is negative for large real parameters, so it is negative on the whole interval $(a_*, +\infty)$. Since $c_* \in (a_*, +\infty)$, the equality above is only possible if $c_* < a_1$. This proves (4.26).

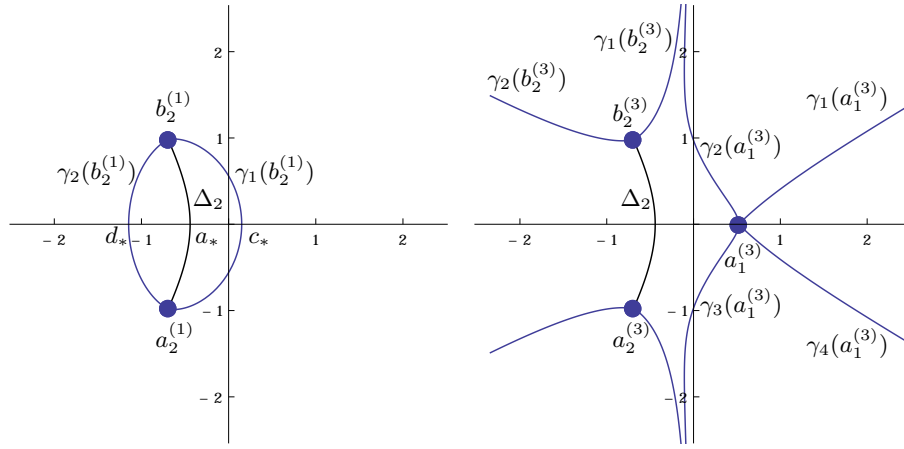


FIGURE 8. $\tau = 0$: blue lines represent the critical graph of ϖ_2 on \mathcal{R}_1 (left) and \mathcal{R}_3 . The black curve is the branch cut Δ_2 connecting \mathcal{R}_1 and \mathcal{R}_3 .

The discussion of the structure of the trajectories $\gamma_1(a_1^{(3)})$, $\gamma_2(a_1^{(3)})$, $\gamma_1(b_2^{(3)})$, $\gamma_2(b_2^{(3)})$ is identical to the analysis of the trajectories $\gamma_j(b_*)$, $\gamma_j(b_2)$, $j = 1, 2$, on \mathcal{R}_2 for ϖ_1 above, so we omit the details.

The global structure of the critical graph on ϖ_2 on both sheets is presented in Figure 8. The basic conclusion is that with the branch cut Δ_2 specified above, the critical graph splits into two sets: a closed Jordan curve on \mathcal{R}_1 , containing $a_2^{(1)}$ and $b_2^{(1)}$, and 4 analytic arcs on \mathcal{R}_3 , starting and ending at $\infty^{(3)}$, each passing through one of the branch points $a_2^{(3)}$, $b_2^{(3)}$ and $a_1^{(3)}$.

4.5.3. Trajectories for $0 < \tau < \frac{1}{12}$. A combination of the general principles **P.2** and **P.6** assures us that the finite critical trajectories for $\tau = 0$ remain finite for small perturbations of τ , and that the behavior of the trajectories described for $\tau = 0$ is preserved for τ small.

Let $\varepsilon > 0$ be sufficiently small. The general principle **P.4** above tells us that if we consider the domains $\Omega_\varepsilon^{(1)}$, (respectively $\Omega_\varepsilon^{(2)}$ and $\Omega_\varepsilon^{(3)}$), swept by trajectories

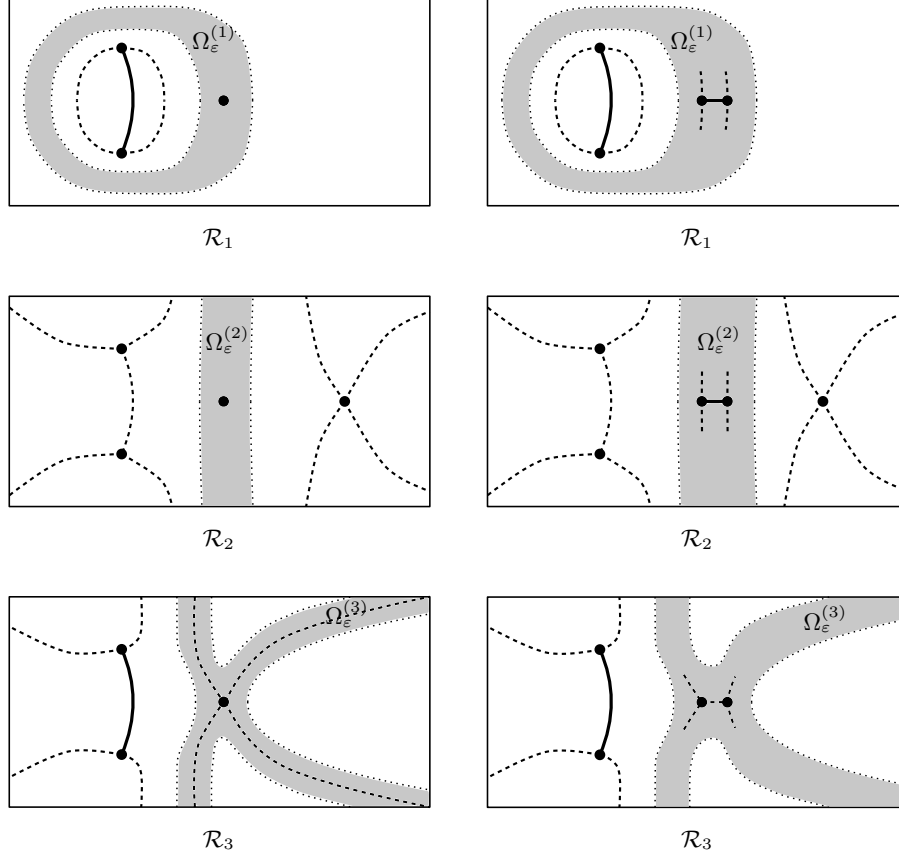


FIGURE 9. Left: critical graph of ϖ for $\tau = 0$, with the domains $\Omega_\varepsilon^{(j)}$, $j = 1, 2, 3$, in gray. Right: local behavior of the critical trajectories for ϖ and $\tau = \delta > 0$, passing through $a_1^{(j)}$, $j = 1, 2, 3$, with the same domains superimposed.

of (4.22) passing through points in the ε -neighborhood of $a_1^{(1)}$ (respectively $a_1^{(2)}$ and $a_1^{(3)}$), then there exists a $\delta > 0$ such that the critical trajectories for ϖ and $0 < \tau < \delta$, passing through $a_1^{(j)}$, $j = 1, 2, 3$, belong to $\Omega_\varepsilon = \Omega_\varepsilon^{(1)} \cup \Omega_\varepsilon^{(2)} \cup \Omega_\varepsilon^{(3)}$. These domains are depicted schematically on Figure 9.

We should keep in mind that the branch cut Δ_2 now is completely specified by Definition 4.5, which is consistent as long as the critical trajectory $\gamma_3(b_2^{(2)})$ joining $a_2^{(2)}$ and $b_2^{(2)}$, *exists and remains on \mathcal{R}_2* for the full range of the parameter τ under consideration. The forthcoming analysis shows that this is the case for $\tau \leq \tau_c$.

According to the general principle **P.4**, there exists a $\delta > 0$ small enough such that for $0 < \tau < \delta$, both $a_1^{(j)}, b_1^{(j)} \in \Omega_\varepsilon^{(j)}$, $j = 1, 2, 3$, and the critical trajectories emerging from $a_1^{(j)}, b_1^{(j)}$ stay in $\Omega_\varepsilon^{(j)}$.

For instance, since for $\tau = 0$ the critical trajectories $\gamma_j(a_1^{(3)})$, $j = 1, \dots, 4$, define three half-plane domains (bounded by $\gamma_1(a_1^{(3)}) \cup \gamma_2(a_1^{(3)})$, $\gamma_3(a_1^{(3)}) \cup \gamma_4(a_1^{(3)})$ and $\gamma_1(a_1^{(3)}) \cup \gamma_4(a_1^{(3)})$, see Figure 8, right), they must be persistent under small perturbation of τ , and either $a_1^{(3)}$ or $b_1^{(3)}$, or both, must belong to their boundaries. Taking into account the structure of $\Omega_\varepsilon^{(3)}$ it is straightforward to conclude that for $0 < \tau < \delta$, the trajectories of ϖ through $a_1^{(3)}$ and $b_1^{(3)}$ are as shown in Figure 10.

We now examine the trajectories from $a_1^{(1)}$ and $b_1^{(1)}$. The following result comes in very handy:

Proposition 4.7. *Let*

$$h(x, y) := \begin{cases} \int_x^y \operatorname{Re}(\xi_{1+}(s) - \xi_3(s)) ds = \int_x^y \operatorname{Re}(\xi_{2+}(s) - \xi_3(s)) ds, & \text{if } x, y \in \Delta_1, \\ \int_x^y \operatorname{Re}(\xi_1(s) - \xi_{2+}(s)) ds = \int_x^y \operatorname{Re}(\xi_1(s) - \xi_{3+}(s)) ds, & \text{if } x, y \in \Delta_3, \end{cases} \quad (4.27)$$

where we integrate along each interval.

If $0 < \tau < 1/12$, then there exists no pair of values $x \neq y$, $x, y \in \Delta_1$, such that $h(x, y) = 0$.

If $1/12 < \tau < 1/4$, and there exists a pair of values $x \neq y$, $x, y \in \Delta_1$ (resp., $x, y \in \Delta_3$) such that $h(x, y) = 0$, then there exists no such pair of values on Δ_3 (resp., on Δ_1).

Furthermore, if $x_1 < y_1$, $x_2 < y_2$ are two such pairs, $(x_1, y_1) \cap (x_2, y_2) \neq \emptyset$.

Notice that h is well defined on Δ_1 and Δ_3 due to the symmetry relations (4.9) and (4.12), and that $\Delta_3 = \emptyset$ for $1/12 < \tau < \tau_c$.

Proof. Assume that there does exist a pair of values $x < y$, $x, y \in \Delta_1$, such that $h(x, y) = 0$ (same analysis is valid for $x, y \in \Delta_3$). By the mean value theorem, there exists a $u \in (x, y)$ such that

$$\operatorname{Re} \xi_{1+}(u) = \operatorname{Re} \xi_{2+}(u) = \xi_3(u),$$

hence,

$$0 = \operatorname{Re} \xi_{1+}(u) + \operatorname{Re} \xi_{2+}(u) + \xi_3(u) = 3\xi_3(u).$$

But

$$D(u) = -\xi_{1+}(u)\xi_{2+}(u)\xi_3(u) = 0,$$

and the assertion follows from Proposition 4.2, keeping in mind that D has exactly one zero on $\Delta_1 \cup \Delta_3$ for $1/12 < \tau < 1/4$. \square

One of the consequences of Proposition 4.7 is that for $0 < \tau < \delta$, the trajectories emanating from $a_1^{(1)}$ and $b_1^{(1)}$ cannot cut Δ_1 and must stay on the sheet \mathcal{R}_1 . Thus, using again the general principle **P.4** we conclude that the trajectories $\gamma_k(a_1^{(1)})$, $\gamma_k(b_1^{(1)})$, $k = 1, 2$, are closed and encircle the cut Δ_2 , see Figure 10.

Similar considerations can be applied to get the behavior for the trajectories emanating from $a_1^{(2)}$, $b_1^{(2)}$, and the final result for τ small is seen in Figure 10. We skip the details.

The outcome of our analysis is that the critical graph of ϖ has the structure showed in Figure 10, at least for $0 < \tau < \delta$. Our next goal is to prove that this is

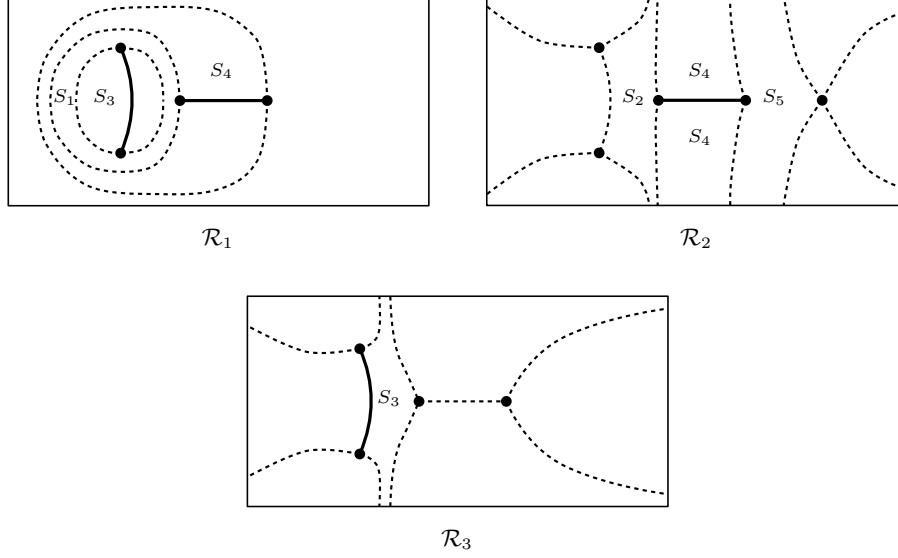


FIGURE 10. Critical graph of ϖ for $0 < \tau < 1/12$, with the strip and ring domains labeled by S_j 's. Notice that some of these domains intersect more than one sheet.

actually valid for $\tau \in (0, 1/12)$. The continuity principle **P.4** yields that this is the case as long as

- (i) No collision of the critical points occur: this is true indeed for $\tau \in (0, 1/12)$, see Section 4.4.
- (ii) No new domains emerge, which amounts to say that the finite trajectories for $0 < \tau < \delta$ remain critical for $0 < \tau < 1/12$: this is assured by a combination of **P.2** and **P.6**.
- (iii) No connected components of the complement of the critical graph “disappear”. More precisely, it means that no width of any strip or ring domain becomes zero. These domains for $0 < \tau < \delta$ are identified on Figure 10: there is one ring domain S_1 and 4 strip domains, S_2, \dots, S_5 . Their widths $\sigma(S_j)$ (see the definition (B.6) in Section B.3) are given by:
 - $\sigma(S_j) = |\omega_j|$, $j = 1, 2, 3$, as defined in (4.19). They do not vanish for $\tau \in (0, 1/12)$, see Figure 4.
 - $\sigma(S_4) = |h(a_1, b_1)|$, with h defined in (4.27), which does not vanish for $\tau \in (0, 1/12)$, see Proposition 4.7.
 - $\sigma(S_5) = \left| \int_{b_1}^{b_*} (\xi_1(s) - \xi_3(s)) ds \right|$, which does not vanish for $\tau \in (0, 1/12)$, see (4.6).

We conclude that the critical graph of ϖ , depicted in Figure 10, is valid for the whole range $0 < \tau < 1/12$. In particular, the critical trajectory $\gamma_1(a_2^{(2)})$ connects the points $a_2^{(2)}$ and $b_2^{(2)}$, which proves Theorem 4.4 for $0 < \tau < 1/12$.

4.5.4. *Trajectories for $\frac{1}{12} < \tau < \tau_1$.* When $\tau = 1/12$, the double point b_* coincides with b_1 , and the strip domain S_5 disappears ($\sigma(S_5) \searrow 0$ as $\tau \nearrow 1/12$), see Figure 11, left. Clearly, this transition has no impact on the structure of trajectories on the

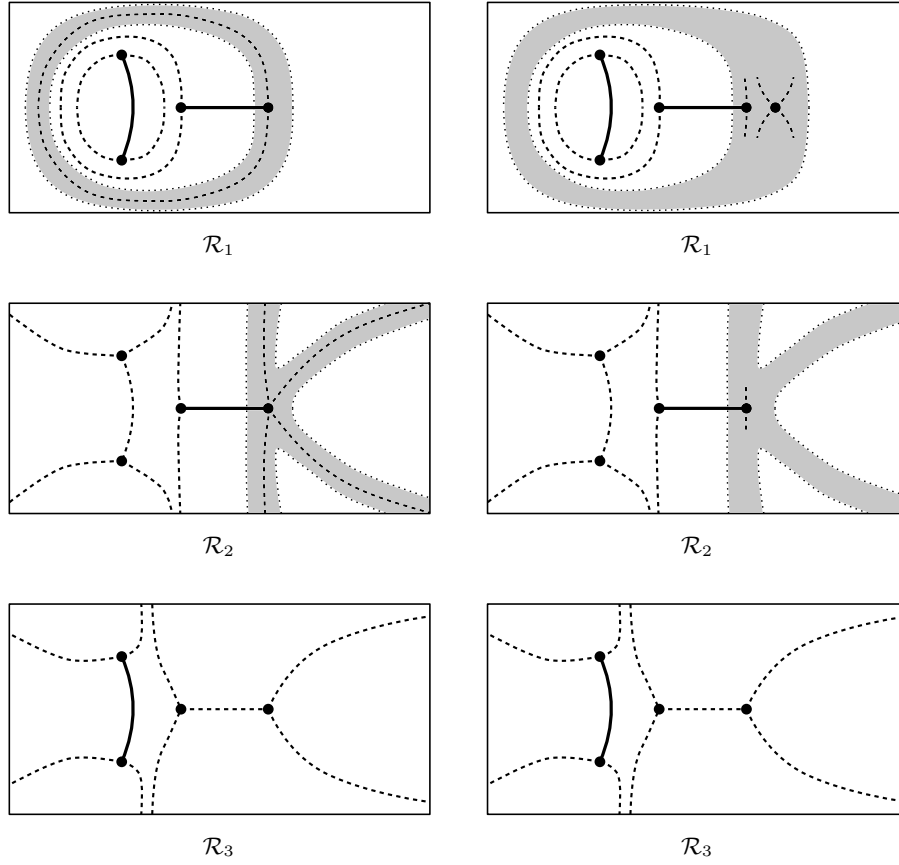


FIGURE 11. Left: critical graph of ϖ for $\tau = 1/12$, with the domains $\Omega_\varepsilon^{(j)}$, $j = 1, 2$, in gray. Right: local behavior of the critical trajectories for ϖ and $\tau = 1/12 + \delta > 0$, passing through $b_1^{(j)}$, $j = 1, 2$, and $b_*^{(1)}$, with the same domains superimposed.

third sheet. Moreover, again a combination of **P.2** and **P.6** assures the finite trajectories for $\tau = 1/12$ remain finite for $1/12 < \tau < 1/12 + \delta$.

Let $\varepsilon > 0$ be sufficiently small. Similarly to what has been done in the previous interval, the general principle **P.4** tells us that if we consider the domains $\Omega_\varepsilon^{(1)}$ and $\Omega_\varepsilon^{(2)}$, swept by trajectories of ϖ passing through points in the ε -neighborhood of $b_1^{(1)}$ and $b_1^{(2)} = b_*^{(2)}$, then there exists a $\delta > 0$ such that the critical trajectories for $1/12 < \tau < 1/12 + \delta$, passing through $b_1^{(j)}$, $j = 1, 2$, and $b_*^{(1)}$, belong to $\Omega_\varepsilon = \Omega_\varepsilon^{(1)} \cup \Omega_\varepsilon^{(2)}$. These domains are also depicted schematically on Figure 11, left.

For $1/12 < \tau < 1/12 + \delta$ we consider the first sheet and the trajectories emanating from $b_1^{(1)}$ and $b_*^{(1)}$; thanks to principle **P.2**, we concentrate on the upper half plane \mathbb{C}_+ (or to be precise, on its pre-image by π on \mathcal{R}_1), namely $\gamma_1(b_*^{(1)})$, $\gamma_2(b_*^{(1)})$,

$\gamma_1(b_1^{(1)})$, see Figure 11, right. These trajectories must stay in Ω_ε , so they have to intersect $\pi^{-1}(\mathbb{R})$ on \mathcal{R}_1 . Let us denote the points of intersection of $\gamma_1(b_*^{(1)})$, $\gamma_2(b_*^{(1)})$, $\gamma_1(b_1^{(1)})$ by $x_1^{(1)}$, $x_2^{(1)}$, $x_3^{(1)}$, respectively. Using the general principles **P.2** and **P.3** we must immediately discard the following possibilities: (i) $x_j \geq b_1$ for some j , (ii) $x_1 \leq a_*^{(1)}$ and $x_2 \leq a_*^{(1)}$ (recall that $a_* = \Delta_2 \cap \mathbb{R}$). Since trajectories cannot intersect, it holds $x_1 < x_2 < x_3$ and we conclude that necessarily $x_2 \in (a_1, b_1)$, and consequently, $x_3 \in (a_1, b_1)$ as well. In particular, $h(x_3, b_1) = 0$, in the notation (4.27).

Since x_1 and x_2 belong to trajectories with a common point $b_*^{(1)}$, the assumption $x_1 \in (a_1, b_1)$ yields that $h(x_1, x_2) = 0$, and since $(x_1, x_2) \cap (x_3, b_1) = \emptyset$, this contradicts Proposition 4.7.

From the considerations above, it follows that $\gamma_1(b_*^{(1)})$ is closed, stays on \mathcal{R}_1 , and intersects $\pi^{-1}(\mathbb{R})$ to the left of Δ_2 , and the trajectories $\gamma_2(b_*^{(1)})$, $\gamma_1(b_1^{(1)})$ intersect the cut Δ_1 and move to the second sheet \mathcal{R}_2 . We keep denoting these points of intersection by x_2, x_3 as before. Clearly, $x_3 > x_2$, and Proposition 4.7 implies that these are the only points of intersection of these trajectories with the interval (a_1, b_1) .

Let us turn to the second sheet, \mathcal{R}_2 and consider $\gamma_1(b_1^{(2)})$: from the structure of $\Omega_\varepsilon^{(2)}$ it is clear that it either diverges to $\infty^{(2)}$, or intersects the branch cut between a_1 and b_1 and moves to \mathcal{R}_1 . If we assume the latter, $\Omega_\varepsilon^{(1)}$ shows that either it will return to \mathcal{R}_2 at a different point in (a_1, b_1) (in contradiction with Proposition 4.7), or it bounds a simply connected domain, which contradicts **P.3**. Hence, $\gamma_1(b_1^{(2)})$ must diverge to $\infty^{(2)}$, and thus $b_1^{(2)}$ lies on the boundary of the half plane domain bounded for $\tau = 1/12$ by $\gamma_1(b_1^{(2)})$ and $\gamma_4(b_1^{(2)})$. In particular, by **P.5**, for $1/12 < \tau < 1/12 + \delta$, $\gamma_1(b_1^{(2)})$ diverges to $\infty^{(2)}$ in the same asymptotic direction given by the angle $\theta_1^{(\infty)}$ from (4.21).

Recall that we concluded that the trajectory $\gamma_1(b_1^{(1)})$ enters \mathcal{R}_2 through the cut (a_1, b_1) at a point x_3 . The only possibility left for it is to go to $\infty^{(2)}$. Applying Theorem B.2 to the ϖ -polygon bounded by the trajectories $\gamma_1(b_1^{(1)})$ and $\gamma_2(b_1^{(2)})$, we get that $\gamma_1(b_1^{(1)})$ goes to $\infty^{(2)}$ in the asymptotic direction given by the angle $\theta_5^{(\infty)}$, and consequently $\gamma_1(b_*^{(1)})$ extends to $\infty^{(2)}$ with angle θ_5 as well.

The outcome of our analysis is that for $1/12 < \tau < 1/12 + \delta$ the critical graph of ϖ has the structure showed in Figure 12. We prove that this is actually valid for $\tau \in (1/12, \tau_1)$. Again, the continuity principle **P.4** yields that this is the case as long as (i) no collision of the critical points occur (this is true indeed for $\tau \in (1/12, 1/4)$, see Section 4.4), (ii) finite critical trajectories for $1/12 < \tau < 1/12 + \delta$ remain finite for $1/12 < \tau < \tau_1$ (assured by a combination of **P.2** and **P.6**); and (iii) no width of any strip and ring domains become zero. These domains for $1/12 < \tau < 1/12 + \delta$ are identified on Figure 12: there is one ring domain S_1 and 4 strip domains, S_2, \dots, S_5 . Their widths $\sigma(S_j)$ are given by:

- $\sigma(S_j) = |\omega_j|$, $j = 1, 2, 3$, as defined in (4.19). They do not vanish for $1/12 < \tau < \tau_1$, see Figure 4, although $\sigma(S_1)$ does vanish for $\tau = \tau_1$.

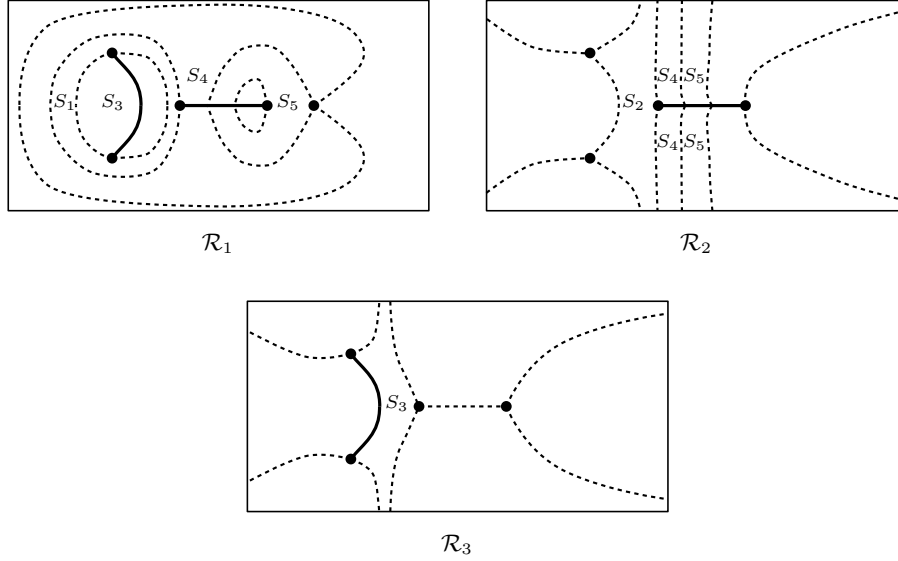


FIGURE 12. Critical graph of ϖ for $1/12 < \tau < \tau_1$, with the strip and ring domains labeled by S_j 's. Notice that some of these domains intersect more than one sheet.

- $\sigma(S_4)$ is given by the absolute value of

$$\begin{aligned} \operatorname{Re} \int_{a_1}^{b_*} (\xi_2(s) - \xi_3(s)) ds &= \operatorname{Re} \int_{a_1}^{b_2} (\xi_2(s) - \xi_3(s)) ds + \operatorname{Re} \int_{b_2}^{b_*} (\xi_2(s) - \xi_3(s)) ds \\ &= -\omega_1 + \omega_4, \end{aligned} \quad (4.28)$$

and $\omega_1 \neq \omega_4$ for $\tau > 1/12$, see Figure 4 in Section 4.3.

- $\sigma(S_5) = \left| \int_{b_1}^{b_*} (\xi_2(s) - \xi_3(s)) ds \right|$, which does not vanish for $\tau > 1/12$, see (4.7).

We conclude that the critical graph of ϖ , depicted in Figure 12, is valid for the whole range $1/12 < \tau < \tau_1$. In particular this yields Theorem 4.4 in the mentioned range of τ .

4.5.5. Trajectories for $\tau_1 < \tau < \tau_c$. At the value $\tau = \tau_1$ the critical trajectory $\gamma_1(b_2^{(1)})$ on the first sheet hits the branch point $a_1^{(1)}$, so that the ring domain S_1 disappears ($\sigma(S_1) \searrow 0$ as $\tau \nearrow \tau_1$), see Figure 12. From the analysis of the behavior of the rest of the widths $\sigma(S_j)$ and the other finite critical trajectories, it follows that this fact does not affect the rest of the strip domains - note that there are no other ring domains. In particular, the trajectories emerging from a_1 , a_2 , and b_2 on the sheets \mathcal{R}_2 , \mathcal{R}_3 do not display any phase transition.

The critical graph for $\tau = \tau_1$ is depicted in Figure 13. In accordance with the methodology we have followed so far, we fix an $\varepsilon > 0$ sufficiently small and consider the domain Ω_ε swept by trajectories of ϖ passing through points in the ε -neighborhood of $b_2^{(1)}$ and $a_1^{(1)} = a_*^{(1)}$. Notice that Ω_ε no longer lives on the single sheet, and its boundary now also contains critical trajectories, namely $\gamma_j(a_1^{(2)})$, $\gamma_j(a_2^{(3)})$, $\gamma_j(b_2^{(3)})$, $j = 1, 2$. This is so because, as already mentioned, there is no

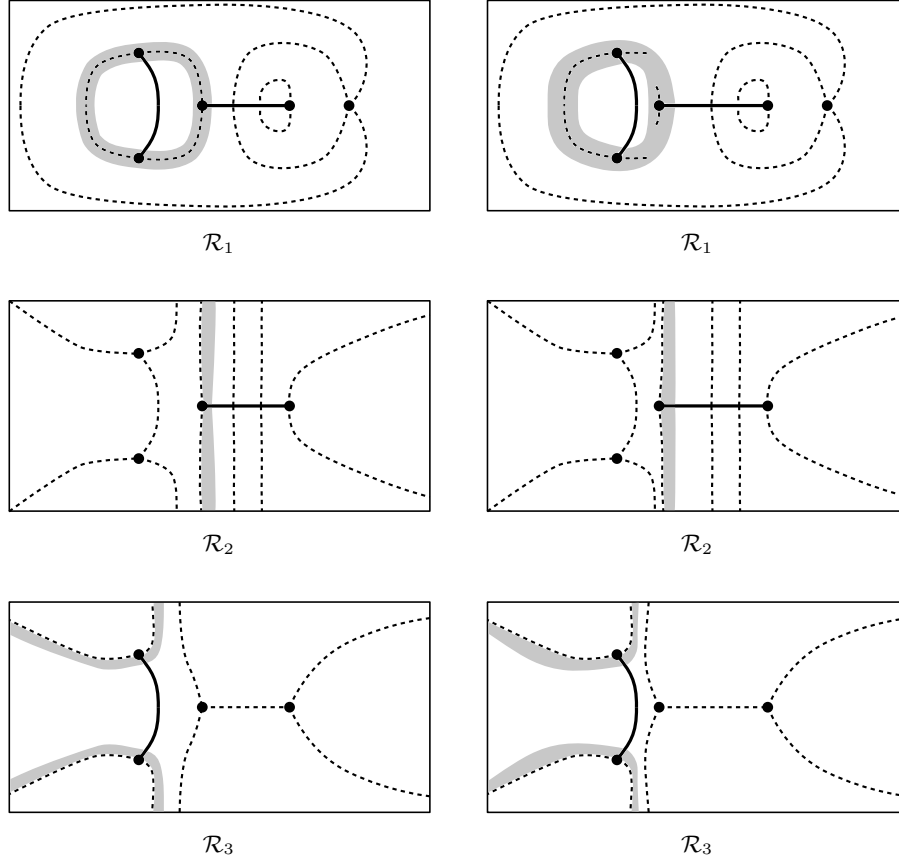


FIGURE 13. Left: critical graph of ϖ for $\tau = \tau_1$, with the domain Ω_ε in gray. Right: local behavior or the critical trajectories for ϖ and $\tau = \tau_1 + \delta > 0$, passing through $b_2^{(1)}$ and $a_1^{(1)}$, with the same domain superimposed.

transition for these trajectories. Observe also that when $\tau \nearrow \tau_1$ the trajectory $\gamma_2(b_2^{(1)}) = \gamma_2(a_2^{(1)})$ does not collide with any critical point other than its endpoints, hence its topology is unchanged under small perturbations of τ around τ_1 . In fact, there exists a $\delta > 0$ such that the critical trajectories for $\tau_1 < \tau < \tau_1 + \delta$, passing through $b_1^{(1)}$, belong to Ω_ε . This domain is also depicted schematically on Figure 13, left.

We now consider the possible behavior of $\gamma_1(b_2^{(1)})$ for $\tau_1 < \tau < \tau_1 + \delta$, having in mind that it cannot leave the shaded region Ω_ε , which shows that either $\gamma_1(b_2^{(1)})$ moves immediately to \mathcal{R}_3 through Δ_2 , or it intersects the (preimage of) the real line near a_1 .

In the first case, $\gamma_1(b_2^{(1)})$ extends to $\infty^{(3)}$ with angle $\theta_3^{(\infty)}$. Thus $\gamma_1(b_2^{(1)}) \cup \gamma_2(b_2^{(3)})$ is the boundary of a ϖ -polygon for which $\kappa = 1$ and $\lambda = 0$, contradicting (B.7).

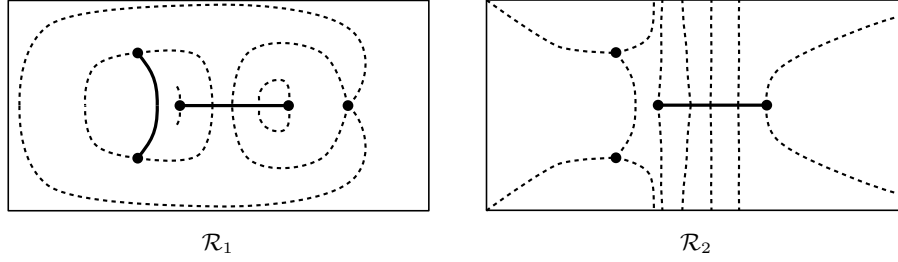


FIGURE 14. Some trajectories of ϖ on the first two sheets, for $\tau_1 < \tau < \tau_c$.

Assume otherwise, so that $\gamma_1(b_2^{(1)})$ intersects the (preimage of the) real line at a point $x^{(1)}$, x close to a_1 . If $x \leq a_1$, we integrate from b_2 to x over $\gamma_1(b_2^{(1)})$ and then from x to a_1 over the real line to get

$$\omega_1 = \operatorname{Re} \int_{b_2}^x (\xi_2(s) - \xi_3(s)) ds + \int_x^{a_1} (\xi_2(s) - \xi_3(s)) ds = \int_x^{a_1} (\xi_2(s) - \xi_3(s)) ds. \quad (4.29)$$

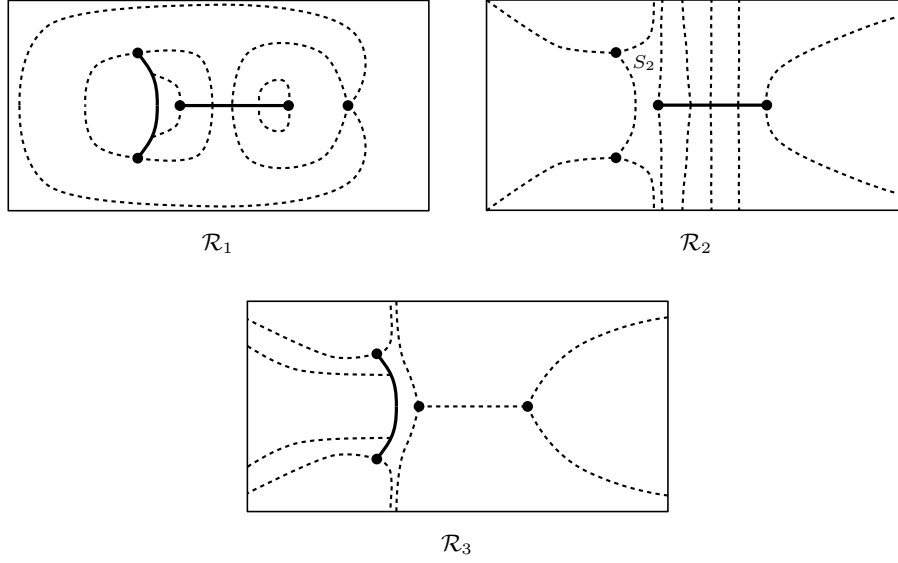
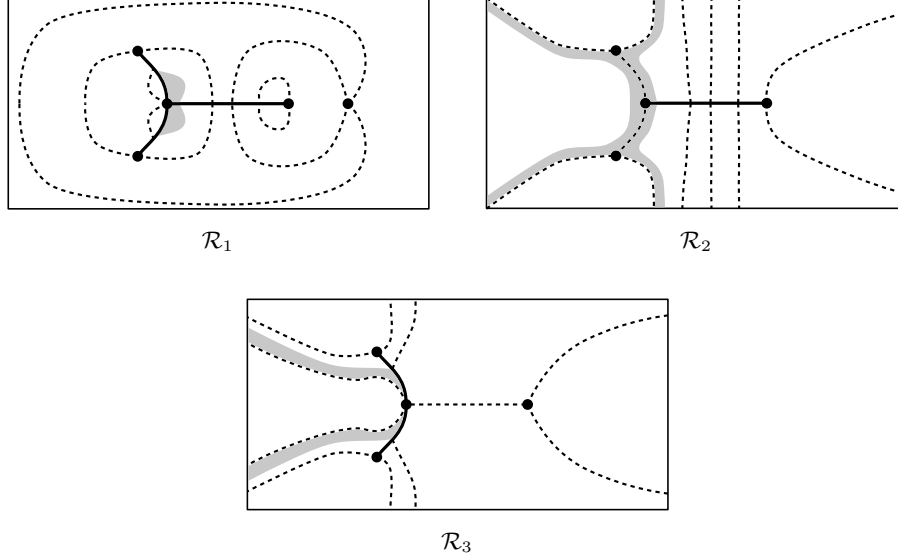
In the range of τ considered, $\omega_1 > 0$, while the last integral is ≤ 0 , see (4.5), which leads us into a contradiction. Hence, $x \in (a_1, b_1)$, so that $\gamma_1(b_2^{(1)})$ moves to the second sheet. Recall that if $x_3^{(1)}$ is the point of intersection of $\gamma_2(b_1^{(1)})$ with Δ_1 , we already have $h(x_3, b_1) = 0$, so by Proposition 4.7, $\gamma_1(b_2^{(1)})$ cannot return to \mathcal{R}_1 ; in consequence, it stays on the second sheet and diverges to $\infty^{(2)}$ in the asymptotic direction corresponding to the angle $\theta_5^{(\infty)}$, see Figure 14.

On the other hand, trajectories $\gamma_j(a_1^{(1)})$, $j = 1, 2$, are “trapped” between $\gamma_1(b_2^{(1)})$ and the branch cut $\Delta_2^{(2)}$, see Figure 14, left. Thus, they cannot stay on \mathcal{R}_1 without violating the general principle **P.3**, so they move to the third sheet and diverge to $\infty^{(3)}$ in the asymptotic directions corresponding to the angles $\theta_3^{(\infty)}$ and $\theta_4^{(\infty)}$, respectively.

The outcome of our analysis is that for $\tau_1 < \tau < \tau_1 + \delta$ the critical graph of ϖ has the structure showed in Figure 15. Same analysis as in the previous section shows that this is actually valid for the whole range $\tau_1 < \tau < \tau_c$, in particular implying Theorem 4.4 for this range of τ .

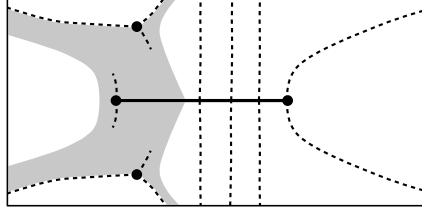
4.5.6. Trajectories for $\tau_c < \tau < \tau_2$. We come to the topologically most important phase transition. According to Definition 4.5, we use the lift of the trajectory joining $a_2^{(2)}$ and $b_2^{(2)}$ on \mathcal{R}_2 as the branch cut connecting the sheets \mathcal{R}_1 and \mathcal{R}_3 . When $\tau = \tau_c$, the strip domain S_2 disappears ($\sigma(S_2) \searrow 0$ as $\tau \nearrow \tau_c$), and this trajectory finally reaches the branch point $a_1^{(2)}$, see Figure 15. Note that this transition corresponds to (1.34). We fix an $\varepsilon > 0$ sufficiently small and consider the domain Ω_ε swept by trajectories of ϖ passing through points in the ε -neighborhood of $a_1^{(2)}$, $a_2^{(2)}$, and $b_2^{(2)}$. The critical graph for $\tau = \tau_c$ along with Ω_ε is displayed in Figure 16.

For $\tau_c < \tau < \tau_c + \delta$ we examine first the second sheet and the trajectories emanating from $a_1^{(2)}$, $a_2^{(2)}$, and $b_2^{(2)}$. For a sufficiently small $\delta > 0$, these trajectories must stay in Ω_ε , see Figure 17.

FIGURE 15. Global structure of trajectories for $\tau_1 < \tau < \tau_c$.FIGURE 16. Critical graph of ϖ for $\tau = \tau_c$, with the domain Ω_ε in gray, consisting of the trajectories for ϖ passing through the ε -neighborhood of $a_1^{(2)}$, $a_2^{(2)}$, and $b_2^{(2)}$.

Analyzing $\gamma_1(a_1^{(2)})$ we have to discard the following possibilities:

- $\gamma_1(a_1^{(2)})$ cannot intersect (the preimage by π of) the real line to the left of a_1 without violating the general principle **P.3**.

FIGURE 17. The domain Ω_ε on \mathcal{R}_2 for $\tau_c < \tau < \tau_c + \delta$.

- it cannot intersect (the preimage by π of) the real line to on the cut (a_1, b_1) either: otherwise the equation $h(x, a_1) = 0$ has a solution in $x \in (a_1, b_1)$, along with the identity $h(x_3, b_1) = 0$, where $x_3^{(1)}$ is the point of intersection of $\gamma_2(b_1^{(1)})$ with Δ_1 , and since $(a_1, x) \cap (x_3, b_1) = \emptyset$, this would contradict Proposition 4.7.
- $\gamma_1(a_1^{(2)}) \neq \gamma_3(b_2^{(2)})$, because $\omega_2 \neq 0$ for $\tau > \tau_c$, see Figure 4.
- $\gamma_1(a_1^{(2)})$ cannot diverge to $\infty^{(2)}$ in the asymptotic direction $\theta_2^{(\infty)}$. Indeed, otherwise either $\gamma_3(b_2^{(2)})$ also diverges to $\infty^{(2)}$ in the same direction, or it intersects the real axis to the left of $a_1^{(2)}$. In the former case, we get a ϖ -polygon for which $\kappa = 1$ and $\lambda = 0$, contradicting (B.7), and in the latter one we proceed as in (4.29) (with ω_1 replaced by ω_2) to get a contradiction.

The only possibility left for $\gamma_1(a_1^{(2)})$ is to diverge to $\infty^{(2)}$ in the asymptotic direction $\theta_3^{(\infty)}$. Since it was already observed that $\gamma_3(b_1^{(2)})$ cannot diverge to $\infty^{(2)}$, it must intersect $\pi^{(-1)}(\mathbb{R})$ to the right of a_1 .

The outcome of our analysis on the second sheet, as well as the region Ω_ε on the remaining sheets, is displayed in Figure 18. The cut Δ_2 is chosen in such a way that its projection on \mathcal{R}_2 coincides with $(\gamma_3(b_2^{(2)}) \cup \gamma_1(a_2^{(2)})) \cap \mathcal{R}_2$.

What is left is to describe the behavior of $\gamma_3(b_2^{(2)})$ on the rest of the sheets. We already saw that this trajectory has to move to \mathcal{R}_1 through the cut Δ_1 . It cannot intersect Δ_1 again (see Proposition 4.7), and it must stay in the region Ω_ε displayed in Figure 18. Hence, the only possibility is that $\gamma_3(b_2^{(2)})$ intersects the cut Δ_2 , moves to the sheet \mathcal{R}_3 and diverges to $\infty^{(3)}$ in the asymptotic direction $\theta_4^{(\infty)}$.

The critical graph of ϖ for $\tau_c < \tau < \tau_c + \delta$ has the structure showed in Figure 19, and we now prove that this is actually valid for $\tau \in (\tau_c, \tau_2)$. Again, the continuity principle **P.4** yields that this is the case as long as (i) no collision of the critical points occur (this is true indeed for $\tau \in (\tau_c, \tau_2)$, see Section 4.4); (ii) finite trajectories for $\tau_c < \tau < \tau_c + \delta$ remain finite for $\tau \in (\tau_c, \tau_2)$ (assured by a combination of **P.2** and **P.6**) and (iii) no width of any strip domains become zero. These domains for $\tau_c < \tau < \tau_c + \delta$ are identified on Figure 19; observe that there are now 7 strip domains S_j , and no ring domains. Let us compute their widths $\sigma(S_j)$.

Lemma 4.8. *For $\tau_c < \tau < \tau_2$,*

$$\sigma(S_1) = \sigma(S_2) = \sigma(S_3) = |\omega_1(\tau)|,$$

with $\omega(\tau)$ defined in (4.19).

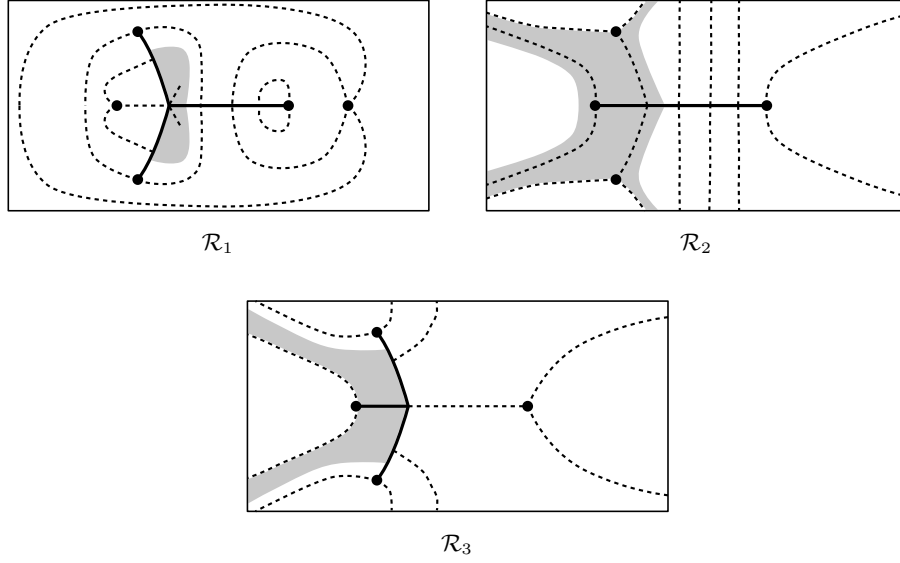


FIGURE 18. Part of the critical graph of ϖ for $\tau_c < \tau < \tau_c + \delta$, with the domain Ω_ε in gray, consisting of the trajectories for ϖ passing through the ε -neighborhood of $a_1^{(2)}$, $a_2^{(2)}$ and $b_2^{(2)}$.

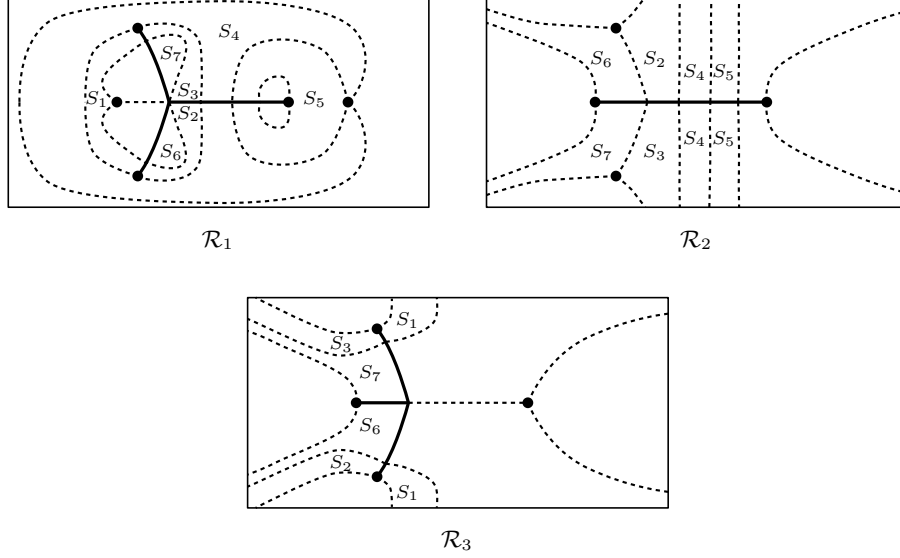


FIGURE 19. Critical graph of ϖ for $\tau_c < \tau < \tau_2$, with the strip domains labeled by S_j 's.

Proof. The fact that $\sigma(S_1) = |\omega_1|$ is the straightforward consequence of the definition of ω_1 . Also $\sigma(S_2) = \sigma(S_3)$ by the symmetry under conjugation. So, it remains to compute the width $\sigma(S_2)$, given by the absolute value of the real part of the

integral

$$\begin{aligned} \int_{b_2^{(2)}}^{a_2^{(1)}=a_2^{(3)}} \sqrt{-\varpi} &= \int_{b_2^{(2)}}^{a_1^{(2)}=a_1^{(3)}} Q(s)ds + \int_{a_1^{(3)}}^{a_2^{(3)}} Q(s)ds \\ &= \int_{b_2}^{a_1} (\xi_1(s) - \xi_3(s))ds + \int_{a_1}^{a_2} (\xi_1(s) - \xi_2(s))ds. \end{aligned}$$

Symmetry under conjugation tells us

$$\operatorname{Re} \int_{a_1}^{a_2} (\xi_1(s) - \xi_2(s))ds = \operatorname{Re} \int_{a_1}^{b_2} (\xi_1(s) - \xi_2(s))ds,$$

and using it in the previous identity, we get

$$\begin{aligned} \sigma(S_2) &= \left| \operatorname{Re} \int_{b_2}^{a_1} (\xi_1(s) - \xi_3(s))ds - \operatorname{Re} \int_{b_2}^{a_1} (\xi_1(s) - \xi_2(s))ds \right| \\ &= \left| \operatorname{Re} \int_{b_2}^{a_1} (\xi_3(s) - \xi_2(s))ds \right| = |\omega_1|. \end{aligned}$$

□

Lemma 4.8 implies that $\sigma(S_j) \neq 0$, $j = 1, 2, 3$, for $\tau_c \leq \tau < \tau_2$, although $\sigma(S_j) \searrow 0$ as $\tau \nearrow \tau_2$, $j = 1, 2, 3$, see Figure 4.

Regarding the rest of the strip domains,

- $\sigma(S_4) = |\omega_4| \neq 0$ for $\tau_c \leq \tau \leq \tau_2$, see again Figure 4.
- $\sigma(S_5) = \left| \int_{b_1}^{b_*} (\xi_2(s) - \xi_3(s))ds \right|$, which does not vanish for $\tau > 1/12$, see (4.7).
- From the symmetry by complex conjugation, $\sigma(S_6) = \sigma(S_7)$, and the structure of S_6 on \mathcal{R}_2 shows that $\sigma(S_6) = |\omega_2|$, so that $\sigma(S_6) = \sigma(S_7) \neq 0$ for $\tau_c < \tau \leq \tau_2$.

Consequently, the critical graph displayed in Figure 19 is valid for $\tau_c < \tau < \tau_2$. In particular it implies Theorem 4.4 for $\tau_c \in (\tau_c, \tau_2)$.

4.5.7. Trajectories for $\tau_2 < \tau < 1/4$. Lemma 4.8 shows that at $\tau = \tau_2$, the strip domains S_1 , S_2 and S_3 displayed in Figure 19 disappear *simultaneously*, which happens because at that moment the branch point $a_1^{(1)}$ hits the critical trajectory $\gamma_1(a_2^{(1)}) = \gamma_2(b_2^{(1)})$ and the point a_* of intersection of Δ_2 with the real line collides with the critical trajectories $\gamma_2(a_2^{(1)})$ and $\gamma_1(b_2^{(1)})$. The resulting critical graph is shown in Figure 20.

As before, we fix an $\varepsilon > 0$ sufficiently small and consider the domain Ω_ε swept by trajectories of ϖ passing through points in the ε -neighborhood of $a_1^{(1)}$, $a_2^{(1)}$, $b_2^{(1)}$ and $a_*^{(j)}$. Observe that for a small perturbation of $\tau = \tau_2$, $\gamma_3(b_2^{(2)}) \cap \mathcal{R}_2$ does not coalesce with critical points other than the starting point $b_2^{(2)}$, hence this arc of critical trajectory displays the same structure as for $\tau = \tau_2$. In particular, the cut Δ_2 is well defined as the projection of $(\gamma_3(b_2^{(2)}) \cup \gamma_1(a_2^{(2)})) \cap \mathcal{R}_2$ on the other sheets. Thus, there exists a $\delta > 0$ such that the critical trajectories for $\tau_2 < \tau < \tau_2 + \delta$, passing through the above mentioned critical points, belong to Ω_ε . This domain is also depicted schematically on Figure 20, left.

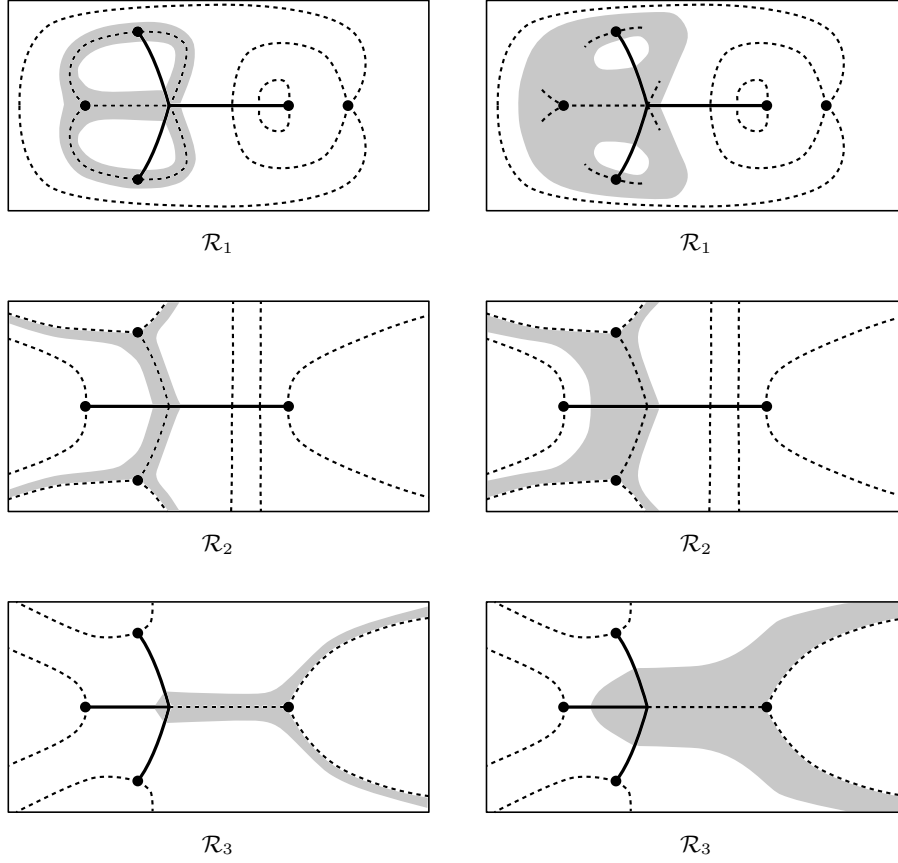


FIGURE 20. Left: critical graph of ϖ for $\tau = \tau_2$, with the domain Ω_ε in gray. Right: local behavior of the critical trajectories for ϖ and $\tau = \tau_2 + \delta > 0$, passing through $a_1^{(1)}$, $a_2^{(1)}$, $b_2^{(1)}$ and a_* , with the same domain superimposed.

Lemma 4.9. *Let t be a point on the part of the curve $\Delta_2^{(1)}$ joining $b_2^{(1)}$ with $a_*^{(1)}$, and γ a Jordan curve on \mathcal{R}_1 connecting the boundary values t_\pm on $\Delta_2^{(1)}$ and containing the only branch point $b_2^{(1)}$ inside. Then*

$$\operatorname{Re} \int_{\gamma} Q(s) ds = \operatorname{Re} \int_{\gamma} (\xi_{2+}(s) - \xi_{3+}(s)) ds = 0.$$

Proof. We can deform γ to the cut Δ_2 . If we denote by $\Delta_2(t)$ the arc of Δ_2 from b_2 to t , we get

$$\begin{aligned} \operatorname{Re} \int_{\gamma} Q(s) ds &= \operatorname{Re} \int_{\Delta_2(t)} (\xi_{2+}(s) - \xi_{3+}(s)) ds - \operatorname{Re} \int_{\Delta_2(t)} (\xi_{2-}(s) - \xi_{3-}(s)) ds \\ &= \operatorname{Re} \int_{\Delta_2(t)} (\xi_{3-}(s) - \xi_{3+}(s)) ds = \operatorname{Re} \int_{\Delta_2(t)} (\xi_{1+}(s) - \xi_{3+}(s)) ds \\ &= \operatorname{Re} \int_{\Delta_2^{(2)}(t)} Q_+ ds = 0, \end{aligned}$$

where we have used the jump condition (4.11) and the fact that $\Delta_2^{(2)}$ is an arc of trajectory of ϖ . \square

Now we describe the critical trajectories for $\tau_2 < \tau < \tau_2 + \delta$, starting with $\gamma_2(a_1^{(1)})$. Notice first that in this case $\omega_1 \neq 0$, so that $\gamma_2(a_1^{(1)})$ cannot contain $b_2^{(1)}$. Furthermore, if $\gamma_2(a_1^{(1)})$ intersects the cut Δ_2 (and thus diverges to $\infty^{(2)}$ in the asymptotic direction $\theta_1^{(\infty)}$), then $\gamma_1(a_1^{(1)}) \cup \gamma_2(a_1^{(1)}) \cup \gamma_1(b_1^{(3)})$ determines a ϖ -polygon for which $\kappa = 1$ and $\lambda = 0$, in a contradiction with (B.7).

Keeping in mind that $\gamma_2(a_1^{(1)})$ must belong to Ω_ε we conclude that it has to intersect the real axis in one of the sheets. Let us denote by x the point of the first intersection of this trajectory with $\pi^{-1}(\mathbb{R})$. Again, we discard some cases:

- x cannot be on \mathcal{R}_1 to the left of a_1 : this yields (by the general principle **P.2**) that $\gamma_2(a_1^{(1)}) = \gamma_3(a_1^{(1)})$, violating **P.3**.
- x cannot be on \mathcal{R}_1 to the right of a_* .

Indeed, if $x > a_*$, we form a curve γ given by the union of three pieces: γ_1 is the arc of $\gamma_1(a_1^{(1)})$ from $a_*^{(1)}$ to $a_1^{(1)}$, γ_2 is the arc of $\gamma_2(a_1^{(1)})$ from $a_1^{(1)}$ to x , and γ_3 is the interval from x to $a_*^{(1)}$. This curve satisfies the assumptions of Lemma 4.9, so that

$$\operatorname{Re} \int_{\gamma} (\xi_{2+}(s) - \xi_{3+}(s)) ds = 0.$$

Since

$$\operatorname{Re} \int_{\gamma_1 \cup \gamma_2} (\xi_{2+}(s) - \xi_{3+}(s)) ds = 0$$

by the definition of trajectories, we conclude that

$$\operatorname{Re} \int_{a_*}^x (\xi_{2+}(s) - \xi_{3+}(s)) ds = 0,$$

that is, $h(a_*, x) = 0$, in the notation (4.27). But if $x_3^{(1)}$ is the point of intersection of $\gamma_2(b_1^{(1)})$ with Δ_1 , we also have $h(x_3, b_1) = 0$, and since $(a_*, x) \cap (x_3, b_1) = \emptyset$, this would contradict Proposition 4.7.

- x cannot be on \mathcal{R}_2 to the left of a_* (in other words, x cannot belong to Δ_3 , which could occur if the trajectory $\gamma_2(a_1^{(1)})$ had slipped to the second sheet through the cut Δ_2 before hitting the real line): this yields

$$\operatorname{Re} \int_x^{a_*} (\xi_{1+}(s) - \xi_{3+}(s)) ds = 0,$$

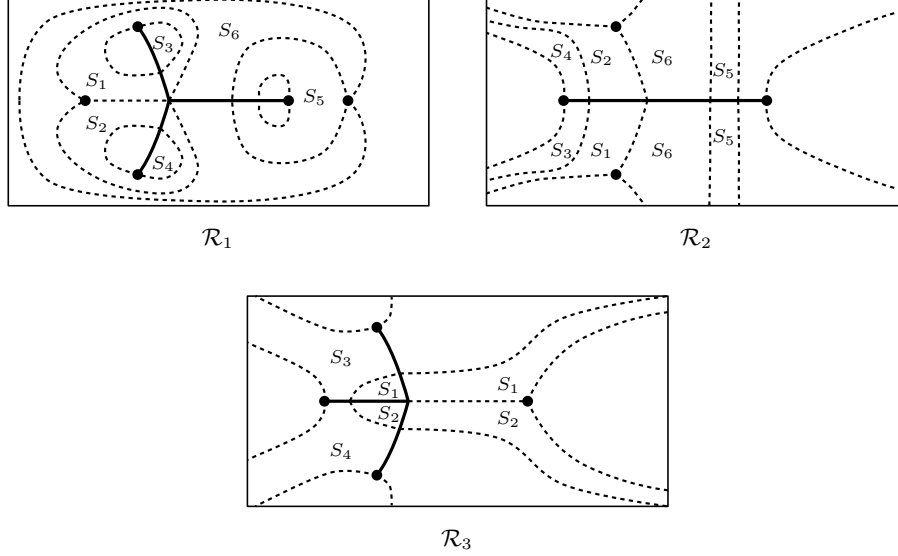


FIGURE 21. Critical graph of ϖ for $\tau_2 < \tau < 1/4$, with the strip domains labeled by S_j 's.

that is, $h(x, a_*) = 0$, in the notation (4.27), with $x, a_* \in \Delta_3$. But if $x_3^{(1)}$ is the point of intersection of $\gamma_2(b_1^{(1)})$ with Δ_1 , we also have $h(x_3, b_1) = 0$, in contradiction with Proposition 4.7.

The only possibility left is that $\gamma_2(a_1^{(1)})$ hits the real line *precisely* at the point a_* . This means also that $\gamma_2(a_1^{(1)}) = \gamma_1(a_2^{(2)})$, as shown in Figure 21.

Regarding the trajectories $\gamma_1(b_2^{(1)})$ and $\gamma_2(b_2^{(1)})$, they cannot satisfy $\gamma_1(b_2^{(1)}) = \gamma_2(b_2^{(1)})$, and due to principle **P.3** they must belong to Ω_ε . Hence, they have to behave as shown in Figure 21. We skip the details.

The critical graph of ϖ for $\tau_2 < \tau < \tau_2 + \delta$ has the structure showed in Figure 21, and we prove that this is actually valid for $\tau \in (\tau_2, 1/4)$. The continuity principle **P.4** yields that this is the case as long as (i) no collision of the critical points occur (this is true indeed for $\tau \in (\tau_2, 1/4)$, see Section 4.4), (ii) the finite critical trajectories for $\tau_2 < \tau < \tau_2 + \delta$ remain finite in the range $\tau \in (\tau_2, 1/4)$ (certainly true by a combination of **P.2** and **P.6**) and (iii) no width of any strip domains become zero. These domains for $\tau_2 < \tau < \tau_2 + \delta$ are identified on Figure 21; observe that there are now 6 strip domains S_j , and no ring domains. Their widths $\sigma(S_j)$ are:

- $\sigma(S_1) = \sigma(S_2) = |\omega_1| \neq 0$ for $\tau \in (\tau_2, 1/4)$.
- $\sigma(S_3) = \sigma(S_4) = |\omega_3| \neq 0$ for $\tau \in (\tau_2, 1/4)$.
- $\sigma(S_5) = \left| \int_{b_1}^{b_*} (\xi_2(s) - \xi_3(s)) ds \right|$, which does not vanish for $\tau > 1/12$, see (4.7).
- As in (4.28), we get that $\sigma(S_6) = |\omega_4 - \omega_1|$, and $\omega_1 \neq \omega_4$ for $\tau > 1/12$, see Figure 4 in Section 4.3.

We conclude that the structure of the critical graph of ϖ , depicted in Figure 21, is actually valid for the whole range $\tau \in (\tau_2, 1/4)$. This finishes the proof of Theorem 4.4.

Remark 4.10. The attentive reader might notice that in Figure 19, for instance, the critical trajectories $\gamma_1(a_2^{(2)})$ and $\gamma_2(a_1^{(1)})$ intersect the cut Δ_2 in \mathcal{R}_1 on pairs of opposite points t_{\pm} . This phenomenon, which also occurs in Figure 21, is easily explained by Lemma 4.9.

Remark 4.11. It follows from the results of Sections 4.5.3–4.5.7 that the trajectories $\gamma_1(b_2^{(2)})$ and $\gamma_2(b_2^{(2)})$ determine a half plane domain H . In particular, this implies that there is an orthogonal critical trajectory γ_1 emerging from $b_2^{(2)}$ which is entirely contained in H and extends to $\infty^{(2)}$ along the angle $2\pi/3$. Similarly, there is an orthogonal critical trajectory γ_2 (which is the complex conjugate of γ_1) emerging from $a_2^{(2)}$ and extending to $\infty^{(2)}$ along the angle $-2\pi/3$. Then the projected contour

$$\Gamma = \pi(\gamma_1 \cup \gamma_2) \cup \Delta_2$$

satisfies the conditions stated in (3.13)–(3.15).

APPENDIX A. PROOF OF THEOREM 1.3 IN THE GENERAL CASE

The proof of Theorem 1.3 presented in Section 2 is valid when the set $\widehat{S}_\alpha \subset S_\alpha$ of double poles of the coefficient D in (1.13) is empty (see Proposition 2.8). Our goal now is to show that for a vector of measures $\vec{\mu} \in \mathcal{M}_\alpha$ whose components are supported on a finite union of analytic arcs, and such the associated functions ξ_j in (1.17) satisfy (1.13) for a polynomial R and a rational function D , the equality

$$\mathcal{D}_h(\vec{\mu}) = 0 \tag{A.1}$$

is valid for every function $h \in C^2(\mathbb{C})$, without any further restriction on the poles of D . Recall that this fact was established so far for the Cauchy kernels h_z defined in (2.3), and when $\text{supp } h \cap \widehat{S}_\alpha = \emptyset$, see Proposition 2.8.

As a first step, we extend (A.1) to polynomials:

Lemma A.1. *Under the assumptions of Theorem 1.3, $\mathcal{D}_q(\vec{\mu}) = 0$ for every algebraic polynomial q .*

Proof. Fix $\rho > 0$ for which

$$\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3 \subset V_\rho := \{x \in \mathbb{C} \mid |x| < \rho\}.$$

For $x \in V_\rho$ and z sufficiently large, we can expand

$$h_z(x) = - \sum_{j=0}^{\infty} \frac{x^j}{z^{j+1}},$$

which implies the identity

$$p_{n-1}(x, z) - z^{n+1}h_z(x) = x^n + w_z(x), \tag{A.2}$$

where the function

$$w_z(x) = \sum_{j=n+1}^{\infty} \frac{x^j}{z^{j-n}}$$

converges to 0 as $z \rightarrow \infty$ uniformly for $x \in V_\rho$, and $p_{n-1}(x, z)$ is a polynomial of degree $n - 1$ in x , given explicitly by

$$p_{n-1}(x, z) = - \sum_{j=0}^{n-1} x^j z^{n-j}.$$

Set $p_{-1} \equiv 0$. If

$$q(x) = \sum_{n=0}^N a_n x^n$$

is a polynomial of degree N , then (A.2) shows that

$$q(x) = -zq(z)h_z(x) + Q_z(x) + W_z(x), \quad (\text{A.3})$$

where

$$Q_z(x) = \sum_{n=0}^N a_n p_{n-1}(x, z), \quad W_z(x) = - \sum_{n=0}^N a_n w_z(x).$$

Note that Q is a polynomial of degree $N - 1$ in x .

Moreover, W_z and W'_z both converge to 0 as $z \rightarrow \infty$ uniformly for $x \in V_\rho$. Hence, the convergences

$$\frac{W_z(x) - W_z(y)}{x - y} \rightarrow 0, \quad \Phi'(x)W_z(x) \rightarrow 0 \quad \text{as } z \rightarrow \infty$$

hold uniformly for $x, y \in V_\rho$, and from the definition of \mathcal{D}_h in (2.1) we get

$$\mathcal{D}_{W_z}(\vec{\mu}) \rightarrow 0 \quad \text{as } z \rightarrow \infty.$$

The quantity $\mathcal{D}_h(\vec{\mu})$ is linear in h , so (A.2) implies

$$\mathcal{D}_q(\vec{\mu}) = -zq(z)\mathcal{D}_{h_z}(\vec{\mu}) + \mathcal{D}_{Q_z}(\vec{\mu}) + \mathcal{D}_{W_z}(\vec{\mu}) = \mathcal{D}_{Q_z}(\vec{\mu}) + \mathcal{D}_{W_z}(\vec{\mu}),$$

where for the last equality we used Corollary 2.7, and hence from the previous limit we get

$$\mathcal{D}_q(\vec{\mu}) = \lim_{z \rightarrow \infty} \mathcal{D}_{Q_z}(\vec{\mu}).$$

The result now follows easily by induction on the degree N of q . If $N = 0$, then Q_z is identically zero, and the equality above implies that $\mathcal{D}_q(\vec{\mu}) = 0$. Assuming now that $\mathcal{D}_h(\vec{\mu}) = 0$ for every polynomial h of degree at most $N - 1$, we get that $\mathcal{D}_{Q_z}(\vec{\mu}) = 0$ because $Q_z(x)$ is a polynomial of degree at most $N - 1$ in x , and the equality above implies $\mathcal{D}_q(\vec{\mu}) = 0$, concluding the proof. \square

If $p \in \widehat{S}_\alpha$, then as it is discussed in Remark 1.10, the point p belongs to exactly two of the supports of μ_1 , μ_2 and μ_3 , and locally the union of these sets is an analytic arc. That is, there exists an open disk U_p centered at p such that

$$\gamma_p = U_p \cap (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3) \quad (\text{A.4})$$

is an analytic arc passing through p .

Additionally, given any point $z \in (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3) \setminus \widehat{S}_\alpha$, there exists a small disk B_z centered z , disjoint from \widehat{S}_α and such that the set

$$B_z \cap (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3)$$

is a finite union of analytic arcs, which can only intersect at the common point z . In case $z \notin S_\alpha$, this intersection reduces to a single analytic arc. The collection

$$\{B_z\} \cup \{U_p\}$$

constructed above is an open cover of the compact set $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$, from which we extract a finite subcover

$$\{B_j\}_{j=1}^m \cup \{U_p\}_{p \in \widehat{S}_p},$$

where for $j = 1, \dots, m$ we have $U_j = U_z$ for some $z = z_j \in (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3) \setminus \widehat{S}_\alpha$. Set

$$B = \bigcup_{j=1}^m B_j, \quad U = \bigcup_{p \in \widehat{S}_\alpha} U_p.$$

It follows from their construction that these sets satisfy

$$B \cap \widehat{S}_\alpha = \emptyset, \quad \text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3 \subset B \cup U. \quad (\text{A.5})$$

Consider a smooth partition of unity $\{\psi_k\}$ of $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$ subordinated to the open cover $B \cup U$. That is, each function ψ_k is real, belongs to $C_\infty(\mathbb{C})$ and additionally satisfies the following properties.

- $0 \leq \psi_k(z) \leq 1$, for every $z \in \mathbb{C}$.
- For every k , $\text{supp } \psi_k \subset U$ or $\text{supp } \psi_k \subset B$.
- Every $z \in \text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$ belongs to the support of a finite number of functions in the collection $\{\psi_k\}$.
- $\sum_k \psi_k(z) = 1$, for every $z \in \text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$.

Since $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$ is compact, the collection $\{\psi_k\}$ can be assumed to be finite. Moreover, we can refine $\{\psi_k\}$ and assume that $\text{supp } \psi_k \subset U$ whenever $\text{supp } \psi_k \cap \widehat{S}_\alpha \neq \emptyset$. Set

$$\widehat{\psi}(z) = \sum_{\text{supp } \psi_k \cap \widehat{S}_\alpha \neq \emptyset} \psi_k(z), \quad \psi(z) = \sum_k \psi_k(z) - \widehat{\psi}(z).$$

The functions ψ and $\widehat{\psi}$ belong to $C_\infty(\mathbb{C})$, satisfy

$$\text{supp } \psi \cap \widehat{S}_\alpha = \emptyset, \quad \text{supp } \widehat{\psi} \subset U \quad (\text{A.6})$$

and

$$\psi(z) + \widehat{\psi}(z) = 1, \quad z \in \text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3. \quad (\text{A.7})$$

Lemma A.2. *Under the conditions of Theorem 1.3, if q is a polynomial, then $\mathcal{D}_{\widehat{\psi}q}(\vec{\mu}) = 0$, where $\widehat{\psi}$ is the function constructed above.*

Proof. From (A.7), it follows that $q \equiv \psi q + \widehat{\psi}q$ on $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$. From the definition of \mathcal{D}_h in (2.1) we then get

$$\mathcal{D}_{\widehat{\psi}q}(\vec{\mu}) = \mathcal{D}_q(\vec{\mu}) - \mathcal{D}_{\psi q}(\vec{\mu}).$$

Using the first condition in (A.6), we get $\text{supp}(\psi q) \cap \widehat{S}_\alpha = \emptyset$, so Lemma 2.8 gives us $\mathcal{D}_{\psi q}(\vec{\mu}) = 0$. Since q is a polynomial, we learn from Lemma A.1 that $\mathcal{D}_q(\vec{\mu}) = 0$, concluding the proof. \square

We are finally able to prove Theorem 1.3 in its full generality.

Proof of Theorem 1.3. Recall the definition of the arcs $\{\gamma_p\}$, $p \in \widehat{S}_\alpha$, given in (A.4). Each of these arcs is a simple contour on the complex plane, and we can find a

smooth arc $\gamma \subset \mathbb{C}$ for which $\cup \gamma_p \subset \gamma$ and $\mathbb{C} \setminus \gamma$ is connected. In particular, the second condition in (A.6) implies

$$\text{supp}(\widehat{\psi}) \cap (\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3) \subset \bigcup_{p \in \widehat{S}_\alpha} \gamma_p \subset \gamma. \quad (\text{A.8})$$

Consider a parametrization $\gamma : [0, 1] \rightarrow \mathbb{C}$ of γ by arc length, set $a = \gamma(0)$ and define a continuous function $g : \gamma \rightarrow \mathbb{C}$ by $g(z) = \gamma'(\gamma^{-1}(z))$. That is, $g(z)$ is a unit vector tangent to γ at the point z , varying continuously with z . Given $h \in C^2(\mathbb{C})$, define

$$G : \gamma \rightarrow \mathbb{C}, \quad G(z) = \frac{1}{g(z)} \left(\frac{\partial h}{\partial x}(z) \text{Re } g(z) + \frac{\partial h}{\partial y}(z) \text{Im } g(z) \right)$$

γ is a simple smooth arc, so it has empty interior. Since $\mathbb{C} \setminus \gamma$ is connected and G is continuous, Mergelyan's Theorem tells us that there exists a sequence of polynomials (p_n) converging to G uniformly on γ . In particular, for $\gamma(t) = z$ this implies that the convergence

$$p_n(\gamma(t))\gamma'(t) = p_n(z)g(z) \rightarrow \frac{\partial h}{\partial x}(z) \text{Re } g(z) + \frac{\partial h}{\partial y}(z) \text{Im } g(z) = \frac{d}{dt}(h(\gamma(t))) \Big|_{t=\gamma^{-1}(z)}$$

holds uniformly on γ , and as a consequence the sequence of polynomials

$$q_n(z) = \int_a^z p_n(s)ds + h(a)$$

converges to

$$\int_0^t \frac{d}{du}(h(\gamma(u)))du + h(a) = h(z)$$

uniformly for $z \in \gamma$.

In summary, we constructed a sequence of polynomials (q_n) converging uniformly to h on γ , and for which the sequence of derivatives $(q'_n) = (p_n)$ converges to G uniformly on γ ; in particular there exists $M > 0$, independent of n , such that

$$|q'_n(x)| \leq M, \quad x \in \gamma, \quad n \geq 1. \quad (\text{A.9})$$

Hence, the convergence

$$\widehat{\psi}(x)q_n(x) \rightarrow \widehat{\psi}(x)h(x) \quad (\text{A.10})$$

holds true uniformly along γ . Due to (A.8), this is enough to conclude that the convergence above holds uniformly on $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$, so that

$$\int \Phi'_j \widehat{\psi} q_n d\mu_j \rightarrow \int \Phi'_j \widehat{\psi} h d\mu_j, \quad j = 1, 2, 3. \quad (\text{A.11})$$

The measures μ_1, μ_2 and μ_3 do not have point masses, so the diagonal $\{x = y\}$ has zero $\mu_j \times \mu_k$ measure. Thus the limit (A.10) also implies that the convergence

$$\frac{\widehat{\psi}(x)q_n(x) - \widehat{\psi}(y)q_n(y)}{x - y} \rightarrow \frac{\widehat{\psi}(x)h(x) - \widehat{\psi}(y)h(y)}{x - y} \quad (\text{A.12})$$

holds true pointwise $(\mu_j \times \mu_k)$ -a.e.

Since the arc γ is connected, we also know that

$$\left| \frac{\widehat{\psi}(x)q_n(x) - \widehat{\psi}(y)q_n(y)}{x - y} \right| \leq \sup_{s \in [0, 1]} \left| \frac{d}{dt} \left((\widehat{\psi} q_n) \circ \gamma \right) (t) \right|_{t=s}$$

whenever $x, y \in \gamma$. In virtue of (A.9) and (A.10), the right-hand side in the inequality above is uniformly bounded in n . Hence the left-hand side of (A.12) is uniformly bounded along γ , and using (A.8) and once again (A.10), we can extend this conclusion to $\text{supp } \mu_1 \cup \text{supp } \mu_2 \cup \text{supp } \mu_3$. Using the Dominated Convergence Theorem and (A.12) we conclude

$$\iint \frac{\widehat{\psi}(x)q_n(x) - \widehat{\psi}(y)q_n(y)}{x - y} d\mu_j(x) d\mu_k(y) \rightarrow \iint \frac{\widehat{\psi}(x)h(x) - \widehat{\psi}(y)h(y)}{x - y} d\mu_j(x) d\mu_k(y),$$

for $j, k = 1, 2, 3$. Combined with (A.11), we finally get

$$\mathcal{D}_{\widehat{\psi}q_n}(\vec{\mu}) \rightarrow \mathcal{D}_{\widehat{\psi}h}(\vec{\mu}).$$

From Lemma A.2, we know that $\mathcal{D}_{\widehat{\psi}q_n}(\vec{\mu}) = 0$ for every n , hence $\mathcal{D}_{\widehat{\psi}h}(\vec{\mu}) = 0$. On the other hand, due to the first condition in (A.6), we have $\text{supp}(\psi h) \cap \widehat{S}_\alpha = \emptyset$, so from Lemma 2.8 we get $\mathcal{D}_{\psi h}(\vec{\mu}) = 0$. Thus,

$$\mathcal{D}_h(\vec{\mu}) = \mathcal{D}_{\widehat{\psi}h}(\vec{\mu}) + \mathcal{D}_{\psi h}(\vec{\mu}) = 0,$$

where for the first equality we used (A.7). Since $h \in C^2(\mathbb{C})$ is arbitrary, Corollary 2.2 gives us that the measure $\vec{\mu}$ is critical, concluding the proof. \square

APPENDIX B. QUADRATIC DIFFERENTIALS

A meromorphic quadratic differential ϖ on a Riemann surface \mathcal{R} is a differential form of type $(2, 0)$, given locally by an expression $f(z)dz^2$, where f is a meromorphic function of a local coordinate z . If $z = z(\zeta)$ is a conformal change of variables, then

$$\widetilde{f}(\zeta)d\zeta^2 = f(z(\zeta))(dz/d\zeta)^2 d\zeta^2$$

represents ϖ in the local coordinates ζ .

In this Appendix, we sketch the minimal background on quadratic differentials used throughout the paper. The general references are the monographs by Strebel [77] and Jenkins [43]; some additional information can be found in [67, 73, 81].

B.1. Critical points and trajectories. The *critical* or *singular points* of $\varpi = f dz^2$ are the zeros and poles of f ; recall that a zero (resp., a pole) of ϖ is a point p where in a local chart sending p to 0 we have $f(z) = z^n \psi(z)$, with $\psi(0) \neq 0$, and with the integer $n \geq 1$ (resp., $n \leq -1$). The value n is the *order* of the critical point p , and is denoted by $\eta(p)$. The rest of the points of \mathcal{R} are called regular, and their order is assumed to be $\eta(p) = 0$.

Critical points of order ≤ -2 (i.e., poles of order 2 and higher) are called *infinite*, and the rest of the critical points are *finite*.

In a neighborhood of any regular point p , the primitive

$$\Upsilon(z) = \int_p^z \sqrt{-\varpi} = \int_p^z \sqrt{-f(s)} ds \quad (\text{B.1})$$

is well defined by specifying the branch of the square root at p and continuing it analytically along the path of integration. Function $\Upsilon(z)$ provides a *distinguished* or a *natural* parameter on \mathcal{R} in a neighborhood of p .

We are mostly interested in the *trajectories* of a quadratic differential ϖ . A Jordan arc $\gamma \subset \mathcal{R}$ is called an *arc of trajectory* of ϖ if it is locally mapped by Υ to

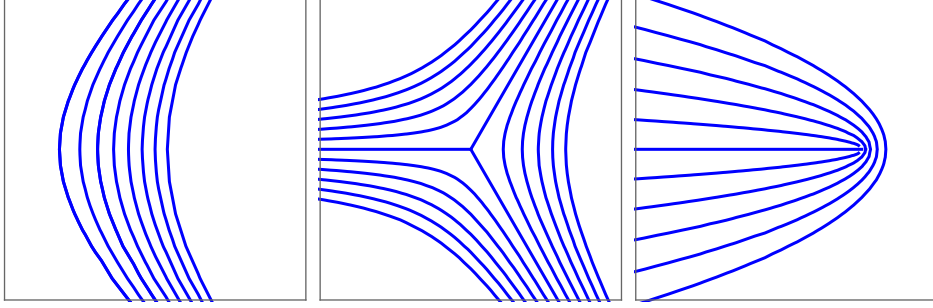


FIGURE 22. Structure of trajectories in a neighborhood of a regular point (left), simple zero (middle) and simple pole (right).

a vertical line. More precisely, this means that for any point $p \in \gamma$, there exists a neighborhood U where the primitive Υ above is well defined and satisfies

$$\operatorname{Re} \Upsilon(z) = \text{const}, \quad z \in \gamma \cap U. \quad (\text{B.2})$$

A maximal arc of trajectory is called a *trajectory* of ϖ .

Analogously, the *orthogonal trajectories* of ϖ are trajectories of $-\varpi$; they can be equivalently defined by replacing “Re” by “Im” in (B.2).

A trajectory γ extending to a finite critical point along at least one of its directions is called *critical*; in the case when it happens in both directions, we call this trajectory *bounded* (also *finite* or *short*), and *unbounded* (or *infinite*) otherwise. Notice that both ends of a short trajectory may coincide, in which case it forms a loop on \mathcal{R} .

A ϖ -*chain* is a connected set on \mathcal{R} made of a finite union of arcs of trajectories or orthogonal trajectories of ϖ . If no curves in a ϖ -chain belong to orthogonal trajectories, we refer to it as a *path of trajectories* of ϖ (or a ϖ -*path*). In this case, in order to avoid the trivial situation, two consecutive arcs of a ϖ -path are required to intersect at a singular point of ϖ .

B.2. The local structure of trajectories. The local behavior of trajectories of a meromorphic quadratic differential ϖ is well understood.

From a point p of order $\eta(p) = n \geq -1$ emanate $n+2$ trajectories, forming equal angles $\frac{2\pi}{n+2}$ at p . This covers also regular points, meaning that through any regular point passes exactly one trajectory, which is locally an analytic arc (see Figure 22).

An infinite critical point p or order $n \leq -3$ has a neighborhood G with the following property: there are $-(n+2)$ asymptotic directions, henceforth called *critical directions*, forming equal angles $\frac{2\pi}{-n-2}$ at p , such that each trajectory entering G stays in G and tends to p in one of the critical directions [43, Theorem 3.3]. If a trajectory is fully contained in G , then it tends to p in two consecutive critical directions (see Figure 23).

At a double pole p there are three possibilities. For $\varpi = f(z)dz^2$, we define the *residue* c of ϖ at $z = p$ to be the residue of $\sqrt{f(z)}$ at $z = p$, which is well defined up to a sign. If $c \in \mathbb{R}$ then there are no trajectories emanating from p and the trajectories near p are closed loops. If $c \in i\mathbb{R}$, then there are trajectories emanating from c in every direction. In the rest of the cases, the trajectories near p converge to p in a spiral form (see Figure 24).

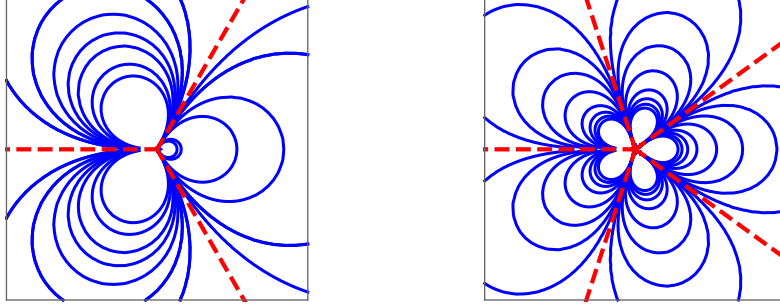


FIGURE 23. Structure of trajectories (solid lines) and critical directions (dashed lines) in a neighborhood of poles of order 3 (left) and 5 (right).

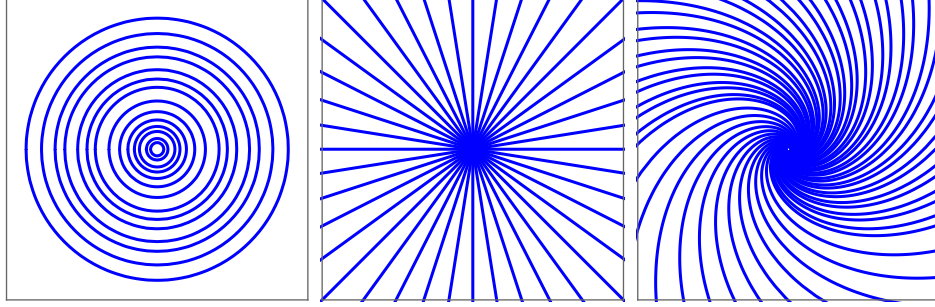


FIGURE 24. Structure of trajectories in a neighborhood of a double pole for $c \in \mathbb{R}$ (left), $c \in i\mathbb{R}$ (middle) and $c \in \mathbb{C} \setminus (\mathbb{R} \cup i\mathbb{R})$ (right).

B.3. Global structure of trajectories. There are three possible behaviors for a trajectory γ in the large,

- (i) γ is a closed curve containing no critical points.
- (ii) γ is an arc connecting two critical points (which may coincide; in this case γ is a closed curve).
- (iii) γ is an arc that has no limit along at least one of its directions.

Trajectories satisfying (ii) are called *short* or *finite*. Trajectories satisfying (iii) are called *recurrent*, and they are usually a major source of troubles when studying the global structure of trajectories of a given quadratic differential. Fortunately (for us) in this paper we deal with quadratic differentials with at most 3 poles on a genus 0 compact Riemann surface, and the absence of recurrent trajectories in this case is assured by Jenkin's Three Poles Theorem [67, Thm. 8.5, page 226].

It is intuitively clear that a non-recurrent and non-closed trajectory γ has two limiting (extremal) values when we travel it in both opposite directions, and that can be distinct or equal. For convenience, we will denote these extremal values by $p(\gamma)$ and $q(\gamma)$. For instance, for a short trajectory γ both $p(\gamma)$ and $q(\gamma)$ are finite critical points.

The set of the critical trajectories of a quadratic differential ϖ (together with their limit points, i.e. the critical points of ϖ) is the *critical graph* of ϖ , denoted

by $\mathcal{G} = \mathcal{G}_\varpi$. According to [43, Theorem 3.5] (see also [77, §10]), the complement of the closure of \mathcal{G}_ϖ in \mathcal{R} consists of a finite number of domains called the *domain configuration* of ϖ . The knowledge of \mathcal{G} (or of the domain configuration of ϖ) is sufficient to fully understand the global structure of trajectories of ϖ , as evidenced by the following theorem:

Theorem B.1 (Basic Structure Theorem [43, Theorem 3.5]). *Let ϖ be a meromorphic quadratic differential on a compact Riemann surface \mathcal{R} . Suppose in addition that ϖ has no recurrent trajectories.*

Then $\mathcal{R} \setminus \mathcal{G}$ decomposes into a finite union of disjoint domains $\cup \mathcal{D}$, each of them bounded by a finite number of critical trajectories. Each domain \mathcal{D} lies into one of the following four classes.

- (i) **Half plane (or end) domain:** *It is swept by trajectories converging to a pole p of order ≥ 3 in its two ends, and along consecutive critical directions. Its boundary consists of a ϖ -path with two unbounded critical trajectories and a finite number of short trajectories. For some choice of the branch of the square root, the natural parameter Υ in (B.1) is a conformal map from \mathcal{D} to a vertical half plane*

$$H = \{w \in \mathbb{C} \mid \operatorname{Re} w > c\},$$

for some $c \in \mathbb{R}$, and it extends continuously to the boundary of \mathcal{D} with the identification $\Upsilon(p) = \infty$.

- (ii) **Strip domain:** *It is swept by trajectories which both ends tend to poles p, q of order ≥ 2 , possibly with $p = q$. The boundary $\partial \mathcal{D} \setminus \{p, q\}$ is a disjoint union of two ϖ -paths, each of them consisting of two unbounded critical trajectories converging to p and q , and possibly a finite number of short trajectories.*

For some real constants $c_1 < c_2$, Υ maps \mathcal{D} conformally to a vertical strip

$$S = \{w \in \mathbb{C} \mid c_1 < \operatorname{Re} w < c_2\}, \quad (\text{B.3})$$

and it extends continuously to the boundary of \mathcal{D} , with appropriate identification of the points p, q with the directions $\pm i\infty$.

- (iii) **Ring domain:** *It is swept by closed trajectories. Its boundary consists of two connected components, where each of them is a closed ϖ -path. For a suitably chosen real constant c and some real numbers $0 < r_1 < r_2$, the function $z \mapsto e^{c\Upsilon(z)}$ maps \mathcal{D} conformally to an annulus*

$$R = \{w \in \mathbb{C} \mid r_1 < |w| < r_2\}. \quad (\text{B.4})$$

and it extends continuously to the boundary of \mathcal{D} .

- (iv) **Circle domain:** *It is swept by closed trajectories and contains exactly one double pole, with purely real residue. Its boundary is a closed ϖ -path.*

For a suitably chosen real constant c and some real number $r > 0$, the function $z \mapsto e^{c\Upsilon(z)}$ is a conformal map from \mathcal{D} to the circle centered at origin and radius r ; it extends continuously to $\partial \mathcal{D}$ and sends the double pole to the origin $w = 0$.

In case ϖ has also recurrent trajectories, a fifth class of domains has to be added to the domain configuration of ϖ ; we refer the reader to [43] for further details.

For a given short trajectory γ connecting two finite critical points p, q (which coincide if γ is closed) we define its *length* by

$$\ell(\gamma) = \left| \int_{\gamma} \sqrt{-\varpi} \right| = \left| \operatorname{Im} \int_{\gamma} \sqrt{-\varpi} \right| > 0, \quad (\text{B.5})$$

The *width of a strip domain* \mathcal{S} is defined as $\sigma(\mathcal{S}) = |c_2 - c_1|$, where the constants c_1, c_2 are as in (B.3). Alternatively, it can be computed as

$$\sigma(\mathcal{S}) = \left| \int_p^q \operatorname{Re} \sqrt{-\varpi} \right|, \quad (\text{B.6})$$

where we integrate along any path in \mathcal{S} connecting two points, p and q , lying on different connected components of $\partial\mathcal{S} \setminus \{\text{poles of } \varpi\}$.

In the same spirit, the *width of a ring domain* \mathcal{D} is defined as

$$\sigma(\mathcal{D}) = \frac{1}{c} \log \frac{r_2}{r_1} = \left| \operatorname{Re} \int_p^q \sqrt{-\varpi} \right|,$$

where r_1, r_2 and c are as in (B.4), and we integrate along any path in \mathcal{D} connecting two points, p and q , lying on different connected components of its boundary.

For any open simply connected domain $\mathcal{D} \subset \mathcal{R}$ bounded by a ϖ -chain (that is called a ϖ -*polygon*) we define two values: the total order of the singular points of ϖ in \mathcal{D} is

$$\lambda(\mathcal{D}) = \sum_{p_j \in \mathcal{D}} \eta(p_j),$$

where the summation is along all the singular points p_j of ϖ in \mathcal{D} , while the contribution from the singular points of ϖ on the boundary is

$$\kappa(\mathcal{D}) = \sum_{p_j \in \partial\mathcal{D}} \beta(p_j),$$

where the summation is along all the corners p_j of $\partial\mathcal{D}$,

$$\beta(p) = 1 - \theta(p) \frac{\eta(p) + 2}{2\pi},$$

and $\theta(p)$ is the inner angle (in radians) at the corner p .

Both values have a simple relation, as shown by the following simple consequence of the argument principle (also known as the Teichmüller lemma, see [77, Theorem 14.1]):

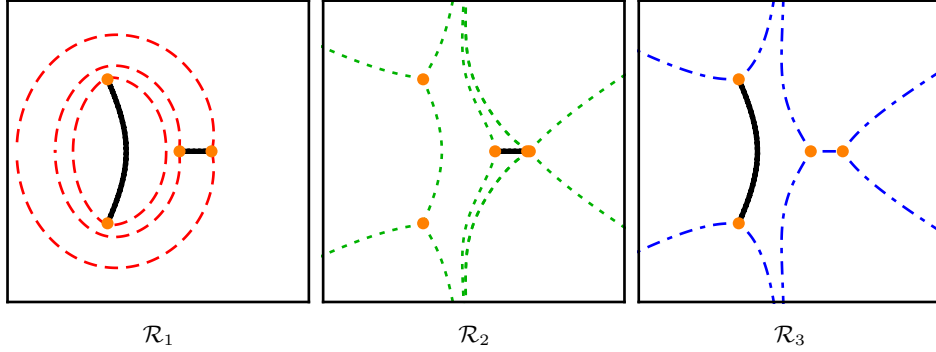
Theorem B.2 (Teichmüller lemma). *If \mathcal{D} is a ϖ -polygon, then*

$$\kappa(\mathcal{D}) = 2 + \lambda(\mathcal{D}). \quad (\text{B.7})$$

A straightforward corollary of this lemma (and also a direct consequence of the maximum principle for holomorphic functions) is the following fact:

Corollary B.3. *If ϖ is analytic (has no poles) in a ϖ -polygon \mathcal{D} , then $\partial\mathcal{D}$ must contain at least one pole of ϖ .*

This corollary is also the basis for the general principle **P.3** in Section 4.5.1.

FIGURE 25. Numerical evaluation of the critical graph for $\tau = 1/25$.

APPENDIX C. NUMERICAL EXPERIMENTS

In this Appendix we look under the hood of the calculation of the functions ω_j in (4.19), as well as of the numerical procedures used to find and plot the trajectories of the quadratic differential ϖ from Section 4.

As it was mentioned, Figure 4 was obtained by means of a numerical evaluation of the integrals defining the functions ω_j in (4.19). For that, we compute the integrals of the form

$$\int_{z_1}^{z_2} \xi_j(s) ds$$

along the line segment joining chosen roots z_1, z_2 of the discriminant in (3.5) by means of the composite trapezoidal rule.

For τ fixed, we compute the points $z_k = z_k(\tau)$, choose a value $m \in \mathbb{N}$, and consider a grid of $m + 1$ equally spaced nodes $\{p_j\}_{j=1}^{m+1}$,

$$p_j = z_1 + \frac{j-1}{m}(z_2 - z_1), \quad j = 1, \dots, m+1,$$

that will be used as the quadrature nodes for the composite trapezoidal rule. At each $z = p_j$ we solve (3.9) numerically, obtaining an (unordered) set of three solutions $\xi_k(p_j)$. Comparing their values at consecutive quadrature points, we collect them into a sequence of (ordered) vectors $\vec{v}_j = (\xi_{\sigma(1)}, \xi_{\sigma(2)}, \xi_{\sigma(3)})^T(p_j)$, $j = 1, \dots, m+1$, where σ is a permutation of $\{1, 2, 3\}$ that does not depend on j . We can then determine the permutation σ using equations (4.8) and (4.14)–(4.17) as boundary conditions.

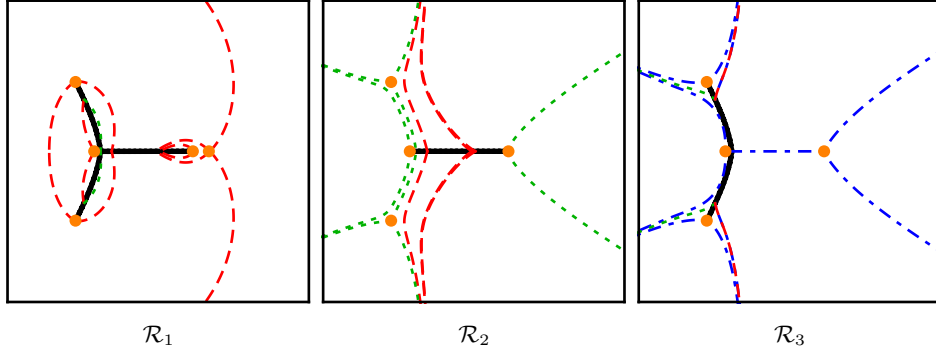
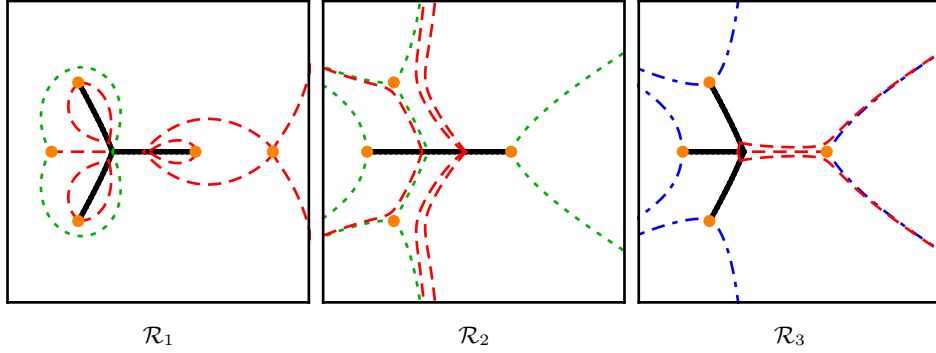
This procedure is repeated for n equally spaced values of τ ,

$$\tau = \tau_k = \frac{k}{4(n+1)}, \quad k = 1, \dots, n,$$

(notice that $\tau_k \in [0, 1/4)$; we stop at $k = n$ in order to avoid the degenerate situation at $\tau = 1/4$ for which some of the endpoints of integration diverge to ∞).

The result of these calculations with $n = 1000$ and $m = 10000$ is plotted in Figure 4.

On the other hand, we also performed numerical experiments that helped us to build the intuition to predict (and confirm) the structure of the critical graphs in Section 4. A sample of such graphs is presented in Figures 25–27.

FIGURE 26. Numerical evaluation of the critical graph for $\tau = 1/5$.FIGURE 27. Numerical evaluation of the critical graph for $\tau = 1/5$.

The numerical procedure used for these pictures is as follows. If $\gamma(t), t \in J \subset \mathbb{R}$, is a parameterization of a trajectory of $\varpi = Q^2(z)dz^2$, then

$$Q^2(\gamma(t))\gamma'(t) = if(t), \quad t \in J,$$

where f is some real valued function. A different choice of f correspond just to a reparameterization of γ . In a natural (arc-length) parametrization of γ the equation above takes the form

$$\gamma'(t) = i \frac{\overline{Q^2(\gamma(t))}}{|Q^2(\gamma(t))|}, \quad t \in J. \quad (\text{C.1})$$

This is an ordinary differential equation for γ that can be solved by standard methods (e.g. the family of Runge-Kutta algorithms).

We should be aware of two aspects when implementing this method. The first one concerns the initial values, that usually are at branch points of (3.9), and thus, at zeros of Q , from where more than one trajectory emanates. This can be resolved by perturbing the initial point along the prescribed direction, which can be obtained from a local analysis of Q at the singular point.

The second aspect concerns the choice of branches of Q^2 in (C.1), where we take advantage of the fact that all the branch points of (3.9) are quadratic. The value of Q^2 is found numerically by solving (3.9), which gives us the three possible values ξ_k . Near a branch point two of these values are close (corresponding to the

two solutions coinciding at this branch point), allowing us to distinguish the third solution. This, in turn, allows to recognize on the the possible branches of Q^2 , but not the other two.

For instance, at $z = b_2$ we will have $\xi_1 \approx \xi_3$ (see (4.17)), which singles out the branch ξ_2 and makes the branch $Q^2 = (\xi_1 - \xi_3)^2$ easily distinguishable. In order to identify the remaining two branches (i.e., $(\xi_2 - \xi_1)^2$ and $(\xi_2 - \xi_3)^2$) we must use further results about the local and global structure of the critical graph of ϖ , established in Section 4.

ACKNOWLEDGEMENTS

Part of this work was carried out during our visit to the Department of Mathematics of the Vanderbilt University, USA, as well as during a stay of the first author as a Visiting Chair Professor at the Department of Mathematics of the Shanghai Jiao Tong University (SJTU), China. We gratefully acknowledge the hospitality and the work environment of the host institutions.

We also thank L. Zhang for explaining his article [35], which was our starting motivation, E. A. Rakhmanov for his interest, and A.B.J. Kuijlaars for a careful reading of the first version of this manuscript and for his useful suggestions. In particular, we give him credit for the Example 1.7.

The first author was partially supported by the Spanish Government together with the European Regional Development Fund (ERDF) under grants MTM2011-28952-C02-01 (from MICINN) and MTM2014-53963-P (from MINECO), by Junta de Andalucía (the Excellence Grant P11-FQM-7276 and the research group FQM-229), and by Campus de Excelencia Internacional del Mar (CEIMAR) of the University of Almería. The second author was supported by FWO Flanders projects G.0641.11 and G.0934.13.

REFERENCES

1. G. Álvarez, L. Martínez-Alonso and E. Medina *Determination of S -curves with applications to the theory of non-Hermitian orthogonal polynomials*, J. Stat. Mech. Theory Exp. (2013), no. 6, P06006, 28.
2. G. Álvarez, L. Martínez-Alonso and E. Medina *Partition functions and the continuum limit in Penner matrix models*, J. Phys. A **47** (2014), no. 31, 315205, 29.
3. A. I. Aptekarev, *Asymptotics of Hermite-Padé approximants for a pair of functions with branch points*, Dokl. Akad. Nauk **422** (2008), no. 4, 443–445.
4. A. I. Aptekarev, V. I. Buslaev, A. Martines-Finkelshtein and S. P. Suetin, *Padé approximants, continued fractions, and orthogonal polynomials*, Uspekhi Mat. Nauk **66** (2011), no. 6(402), 37–122, translation in Russian Math. Surveys **66** (2011), no. 6, 1049–1131.
5. A. I. Aptekarev and A. B. È. Koïlaars, *Hermite-Padé approximations and ensembles of multiple orthogonal polynomials*, Uspekhi Mat. Nauk **66** (2011), no. 6(402), 123–190.
6. A. I. Aptekarev, A. B. J. Kuijlaars and W. Van Assche, *Asymptotics of Hermite-Padé rational approximants for two analytic functions with separated pairs of branch points (case of genus 0)*, Int. Math. Res. Pap. IMRP (2008), Art. ID rpm007, 128.
7. A. I. Aptekarev, D. N. Toulyakov and W. Van Assche, *Hyperelliptic uniformization of algebraic curves of the third order*, J. Comput. Appl. Math. **284** (2015), 38–49.
8. A. I. Aptekarev, W. Van Assche and M. L. Yatsselev, *Hermite-Padé approximants for a pair of Cauchy transforms with overlapping symmetric supports*, preprint Arxiv:math.1505.03993, 2015.
9. M. J. Atia, A. Martínez-Finkelshtein, P. Martínez-González and F. Thabet, *Quadratic differentials and asymptotics of Laguerre polynomials with varying complex parameters*, J. Math. Anal. Appl. **416** (2014), no. 1, 52–80.
10. J. Baik, P. Deift, K. T.-R. McLaughlin, P. Miller and X. Zhou, *Optimal tail estimates for directed last passage site percolation with geometric random variables*, Adv. Theor. Math. Phys. **5** (2001), no. 6, 1207–1250.
11. F. Balogh and M. Bertola, *Regularity of a vector potential problem and its spectral curve*, J. Approx. Theory **161** (2009), no. 1, 353–370.
12. F. Balogh, M. Bertola and T. Bothner, *Hankel determinant approach to generalized Vorob’ev-Yablonski polynomials and their roots*, preprint Arxiv:math.1504.00440, 2015.
13. B. Beckermann, V. Kalyagin, A. C. Matos and F. Wielonsky, *Equilibrium problems for vector potentials with semidefinite interaction matrices and constrained masses*, Constr. Approx. **37** (2013), no. 1, 101–134.
14. T. Bergkvist and H. Rullgård, *On polynomial eigenfunctions for a class of differential operators*, Math. Res. Lett. **9** (2002), no. 2-3, 153–171.
15. M. Bertola, *Boutroux curves with external field: equilibrium measures without a variational problem*, Anal. Math. Phys. **1** (2011), no. 2-3, 167–211.
16. M. Bertola, M. Gekhtman and J. Szmigielski, *Strong asymptotics for Cauchy biorthogonal polynomials with application to the Cauchy two-matrix model*, J. Math. Phys. **54** (2013), no. 4, 043517, 25.
17. M. Bertola and T. Bothner, *Zeros of large degree Vorob’ev-Yablonski polynomials via a Hankel determinant identity*, preprint Arxiv:math.1401.1408, 2014.
18. M. Bertola and A. Tovbis, *Asymptotics of orthogonal polynomials with complex varying quartic weight: global structure, critical point behaviour and the first Painlevé equation*, Constr. Approx. **41** (2015), no. 3, 529–587.
19. J.-E. Björk, J. Borcea and R. Bøgvad, *Subharmonic configurations and algebraic Cauchy transforms of probability measures*, Notions of positivity and the geometry of polynomials, Trends Math., Birkhäuser/Springer Basel AG, Basel, 2011, pp. 39–62.
20. P. M. Bleher and A. Deaño, *Painlevé I double scaling limit in the cubic matrix model*, preprint Arxiv:math.1301.4266, 2013.
21. P. M. Bleher, S. Delvaux and A. B. J. Kuijlaars, *Random matrix model with external source and a constrained vector equilibrium problem*, Comm. Pure Appl. Math. **64** (2011), no. 1, 116–160.
22. P. M. Bleher and A. B. J. Kuijlaars, *Large n limit of Gaussian random matrices with external source. I*, Comm. Math. Phys. **252** (2004), no. 1-3, 43–76.

23. P. M. Bleher and A. B. J. Kuijlaars, *Random matrices with external source and multiple orthogonal polynomials*, Int. Math. Res. Not. (2004), no. 3, 109–129.
24. P. M. Bleher and K. Liechty, *Random matrices and the six-vertex model*, CRM Monograph Series, vol. 32, American Mathematical Society, Providence, RI, 2014.
25. P. M. Bleher and A. Deaño, *Topological expansion in the cubic random matrix model*, International Mathematics Research Notices (2013), no. 12, 2699–2755.
26. P. M. Bleher and A. B. J. Kuijlaars, *Orthogonal polynomials in the normal matrix model with a cubic potential*, Adv. Math. **230** (2012), no. 3, 1272–1321.
27. A. Deaño, D. Huybrechs and A. B. J. Kuijlaars, *Asymptotic zero distribution of complex orthogonal polynomials associated with Gaussian quadrature*, J. Approx. Theory **162** (2010), no. 12, 2202–2224.
28. P. Deift, *Orthogonal polynomials and random matrices: A Riemann-Hilbert approach*, Courant Lecture Notes, no. 3, American Mathematical Society, 2000.
29. P. Deift, T. Kriecherbauer, K. T-R McLaughlin, S. Venakides and X. Zhou, *Uniform asymptotics for polynomials orthogonal with respect to varying exponential weights and applications to universality questions in random matrix theory*, Comm. Pure Appl. Math. **52** (1999), 1335–1425.
30. P. Deift, A. Its and I. Krasovsky, *Asymptotics of Toeplitz, Hankel, and Toeplitz+Hankel determinants with Fisher-Hartwig singularities*, Ann. of Math. (2) **174** (2011), no. 2, 1243–1299.
31. P. Deift, A. Its and I. Krasovsky, *Toeplitz matrices and Toeplitz determinants under the impetus of the Ising model: some history and some recent results*, Comm. Pure Appl. Math. **66** (2013), no. 9, 1360–1438.
32. M. Duits, D. Geudens and A. B. J. Kuijlaars, *A vector equilibrium problem for the two-matrix model in the quartic/quadratic case*, Nonlinearity **24** (2011), no. 3, 951–993.
33. M. Duits and A. B. J. Kuijlaars, *Universality in the two-matrix model: a Riemann-Hilbert steepest-descent analysis*, Comm. Pure Appl. Math. **62** (2009), no. 8, 1076–1153.
34. M. Duits, A. B. J. Kuijlaars and M. Y. Mo, *The Hermitian two matrix model with an even quartic potential*, Mem. Amer. Math. Soc. **217** (2012), no. 1022, v+105.
35. G. Filipuk, W. Van Assche and L. Zhang, *Multiple orthogonal polynomials associated with an exponential cubic weight*, J. Approx. Theory **190** (2015), 1–37.
36. A. S. Fokas, A. Its and A. V. Kitaev, *The isomonodromy approach to matrix models in 2D quantum gravity*, Comm. Math. Phys. **147** (1992), 395–430.
37. J. S. Geronimo, A. B. J. Kuijlaars and W. Van Assche, *Riemann-Hilbert problems for multiple orthogonal polynomials*, Special functions 2000: current perspective and future directions (Tempe, AZ), NATO Sci. Ser. II Math. Phys. Chem., vol. 30, Kluwer Acad. Publ., Dordrecht, 2001, pp. 23–59.
38. A. A. Gonchar and E. A. Rakhmanov, *On the convergence of simultaneous Padé approximants for systems of functions of Markov type*, Trudy Mat. Inst. Steklov. **157** (1981), 31–48, 234, Number theory, mathematical analysis and their applications.
39. A. A. Gonchar and E. A. Rakhmanov, *The equilibrium problem for vector potentials*, Uspekhi Mat. Nauk **40** (1985), no. 4(244), 155–156.
40. A. A. Gonchar and E. A. Rakhmanov, *Equilibrium distributions and the rate of rational approximation of analytic functions*, Mat. Sb. (N.S.) **134(176)** (1987), no. 3, 306–352, 447.
41. A. Hardy and A. B. J. Kuijlaars, *Weakly admissible vector equilibrium problems*, J. Approx. Theory **164** (2012), no. 6, 854–868.
42. T. Holst and B. Shapiro, *On higher Heine-Stieltjes polynomials*, Israel J. Math. **183** (2011), 321–345.
43. J. A. Jenkins, *Univalent functions and conformal mapping*, Ergebnisse der Mathematik und ihrer Grenzgebiete. Neue Folge, Heft 18. Reihe: Moderne Funktionentheorie, Springer-Verlag, Berlin-Göttingen-Heidelberg, 1958.
44. V. A. Kalyagin, *A class of polynomials determined by two orthogonality relations*, Mat. Sb. (N.S.) **110(152)** (1979), no. 4, 609–627.
45. S. Kamvissis, *Comment on “Existence and regularity for an energy maximization problem in two dimensions” [S. Kamvissis and E. A. Rakhmanov, J. Math. Phys. **46**, 083505 (2005)]*, J. Math. Phys. **50** (2009), 104101, 3.
46. S. Kamvissis and E. A. Rakhmanov, *Existence and regularity for an energy maximization problem in two dimensions*, J. Math. Phys. **46** (2005), no. 8, 083505, 24.

47. A. B. J. Kuijlaars, *Multiple orthogonal polynomials in random matrix theory*, Proceedings of the International Congress of Mathematicians. Volume III (New Delhi), Hindustan Book Agency, 2010, pp. 1417–1432.
48. A. B. J. Kuijlaars, A. Martínez-Finkelshtein and F. Wielonsky, *Non-intersecting squared Bessel paths and multiple orthogonal polynomials for modified Bessel weights*, Comm. Math. Phys. **286** (2009), no. 1, 217–275.
49. A. B. J. Kuijlaars and K. T.-R. McLaughlin, *Riemann-Hilbert analysis for Laguerre polynomials with large negative parameter*, Comput. Methods Funct. Theory **1** (2001), no. 1, 205–233.
50. A. B. J. Kuijlaars and K. T.-R. McLaughlin, *Asymptotic zero behavior of Laguerre polynomials with negative parameter*, Constr. Approx. **20** (2004), no. 4, 497–523.
51. A. B. J. Kuijlaars, W. Van Assche and F. Wielonsky, *Quadratic Hermite-Padé approximation to the exponential function: a Riemann-Hilbert approach*, Constr. Approx. **21** (2005), no. 3, 351–412.
52. A. B. J. Kuijlaars and A. López-García, *The normal matrix model with a monomial potential, a vector equilibrium problem, and multiple orthogonal polynomials on a star*, Nonlinearity **28** (2015), no. 2, 347–406.
53. A. B. J. Kuijlaars and G. L. F. Silva, *S-curves in polynomial external fields*, J. Approx. Theory **191** (2015), 1–37.
54. A. B. J. Kuijlaars and A. Tovbis, *The supercritical regime in the normal matrix model with cubic potential*, Adv. Math. **283** (2015), 530–587.
55. G. V. Kuz'mina, *Moduli of families of curves and quadratic differentials*, Proc. Steklov Inst. Math. **139** (1982), 1–231.
56. O. Lehto, *Univalent functions and Teichmüller spaces*, Graduate Texts in Mathematics, vol. 109, Springer-Verlag, New York, 1987.
57. A. Martínez-Finkelshtein, P. Martínez-González and F. Thabet, *Trajectories of quadratic differentials for Jacobi polynomials with complex parameters*, preprint Arxiv:math.1506.03434, 2015.
58. A. Martínez-Finkelshtein and E. A. Rakhmanov, *Critical measures, quadratic differentials, and weak limits of zeros of Stieltjes polynomials*, Comm. Math. Phys. **302** (2011), no. 1, 53–111.
59. A. Martínez-Finkelshtein and G. L. F. Silva, *Critical measures for vector energy: breaking the symmetry for multiple orthogonal polynomials*, in preparation, 2016.
60. A. Martínez-Finkelshtein and E. A. Rakhmanov, *On asymptotic behavior of Heine-Stieltjes and Van Vleck polynomials*, Recent trends in orthogonal polynomials and approximation theory, Contemp. Math., vol. 507, Amer. Math. Soc., Providence, RI, 2010, pp. 209–232.
61. A. Martínez-Finkelshtein, E. A. Rakhmanov and S. P. Suetin, *Heine, Hilbert, Padé, Riemann, and Stieltjes: John Nuttall's work 25 years later*, Recent advances in orthogonal polynomials, special functions, and their applications, Contemp. Math., vol. 578, Amer. Math. Soc., Providence, RI, 2012, pp. 165–193.
62. V. N. Monakhov and L. I. A. Leifman, *Boundary-value problems with free boundaries for elliptic systems of equations*, Translations of mathematical monographs, American Mathematical Society, 1983.
63. E. M. Nikishin and V. N. Sorokin, *Rational approximations and orthogonality*, Translations of Mathematical Monographs, vol. 92, American Mathematical Society, Providence, RI, 1991, Translated from the Russian by R. P. Boas.
64. J. Nuttall, *Padé polynomial asymptotics from a singular integral equation*, Constr. Approx. **6** (1990), no. 2, 157–166.
65. J. Ortega-Cerdà and B. Pridhnani, *The Pólya-Tchebotaröv problem*, Harmonic analysis and partial differential equations, Contemp. Math., vol. 505, Amer. Math. Soc., Providence, RI, 2010, pp. 153–170.
66. E. Perevoznikova and E. A. Rakhmanov, *Variations of the equilibrium energy and S-property of compacta of minimal capacity*, manuscript, 1994.
67. C. Pommerenke, *Univalent functions - With a chapter on quadratic differentials by G. Jensen*, Vandenhoeck & Ruprecht, Göttingen, 1975.
68. E. A. Rakhmanov, *On the asymptotics of Hermite-Padé polynomials for two Markov functions*, Mat. Sb. **202** (2011), no. 1, 133–140.

69. E. A. Rakhmanov, *Orthogonal polynomials and S-curves*, Contemp. Math., vol. 578, Amer. Math. Soc., Providence, RI, 2012.
70. E. A. Rakhmanov, *The asymptotics of Hermite-Padé polynomials for two Markov-type functions*, Sbornik: Mathematics **202** (2011), no. 1, 127.
71. M. Schiffer, *Variation of domain functionals*, Bull. Amer. Math. Soc. **60** (1954), 303–328.
72. M. Schiffer, *A Method of Variation Within the Family of Simple Functions*, Proc. London Math. Soc. **S2-44**, no. 6, 432.
73. A. Y. Solynin, *Quadratic differentials and weighted graphs on compact surfaces*, Analysis and mathematical physics, Trends Math., Birkhäuser, Basel, 2009, pp. 473–505.
74. H. Stahl, *Orthogonal polynomials with complex-valued weight function. I,II*, Constr. Approx. **2** (1986), no. 3, 225–240, 241–251.
75. H. Stahl, *Orthogonal polynomials with respect to complex-valued measures*, Orthogonal polynomials and their applications (Erice, 1990), IMACS Ann. Comput. Appl. Math., vol. 9, Baltzer, Basel, 1991, pp. 139–154.
76. T. J. Stieltjes, *Sur certains polynômes que vérifient une équation différentielle linéaire du second ordre et sur la théorie des fonctions de Lamé*, Acta Math. **6** (1885), 321–326.
77. K. Strebel, *Quadratic differentials*, Ergebnisse der Mathematik und ihrer Grenzgebiete (3) [Results in Mathematics and Related Areas (3)], vol. 5, Springer-Verlag, Berlin, 1984.
78. G. Szegő, *Orthogonal polynomials*, fourth ed., American Mathematical Society, Colloquium Publications, no. 23, American Mathematical Society, 1975.
79. O. Teichmüller, *Extremale quasikonforme Abbildungen und quadratische Differentiale*, Abh. Preuss. Akad. Wiss. Math.-Nat. Kl. **1939** (1940), no. 22, 197.
80. W. Van Assche, *Padé and Hermite-Padé approximation and orthogonality*, Surv. Approx. Theory **2** (2006), 61–91.
81. A. Vasiliev, *Moduli of families of curves for conformal and quasiconformal mappings*, Lecture Notes in Mathematics, vol. 1788, Springer-Verlag, Berlin, 2002.
82. S. Zdravkovska, *Conversation with Vladimir Igorevich Arnol'd (Arnold interviewed by Smilka Zdravkovska)*, The Mathematical Intelligencer **9** (1987), no. 4, 28–32.

(AMF) DEPARTMENT OF MATHEMATICS, UNIVERSITY OF ALMERÍA, ALMERÍA, SPAIN
E-mail address: `andrei@ual.es`

(GS) KU LEUVEN, DEPARTMENT OF MATHEMATICS, CELESTIJNENLAAN 200B BUS 2400, B-3001 LEUVEN, BELGIUM

*CURRENT ADDRESS: DEPARTMENT OF MATHEMATICS, UNIVERSITY OF MICHIGAN, 530 CHURCH STREET, ANN ARBOR, MI 48109, USA
E-mail address: `silvag@umich.edu`